THE UNIVERSITY OF ALBERTA

A LINEAR BOUNDARY COLLOCATION METHOD FOR SOLVING

THE DIRICHLET PROBLEM FOR LAPLACE'S EQUATION

by

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A THESIS

SUBMITTED TO THE FACULTY OF GRADUATE STUDIES

IN PARTIAL FULFILMENT OF THE REQUIREMENTS FOR THE DEGREE

OF MASTER OF SCIENCE

DEPARTMENT OF COMPUTING SCIENCE

EDMONTON, ALBERTA

JULY, 1968

UNIVERSITY OF ALBERTA

FACULTY OF GRADUATE STUDIES

The undersigned certify that they have read, and
recommend to the Faculty of Graduate Studies for acceptance,
a thesis entitled A LINEAR BOUNDARY COLLOCATION METHOD
FOR SOLVING THE DIRICHLET PROBLEM FOR LAPLACE'S EQUATION
submitted by Clifford G. Morgan in partial fulfilment
of the requirements for the degree of Master of Science.

ABSTRACT

This thesis presents a computationally-oriented
linear boundary collocation technique for obtaining a
close approximation to the solution of the Dirichlet
problem for Laplace's equation. Necessary concepts and
algorithms from the theory of linear approximation are
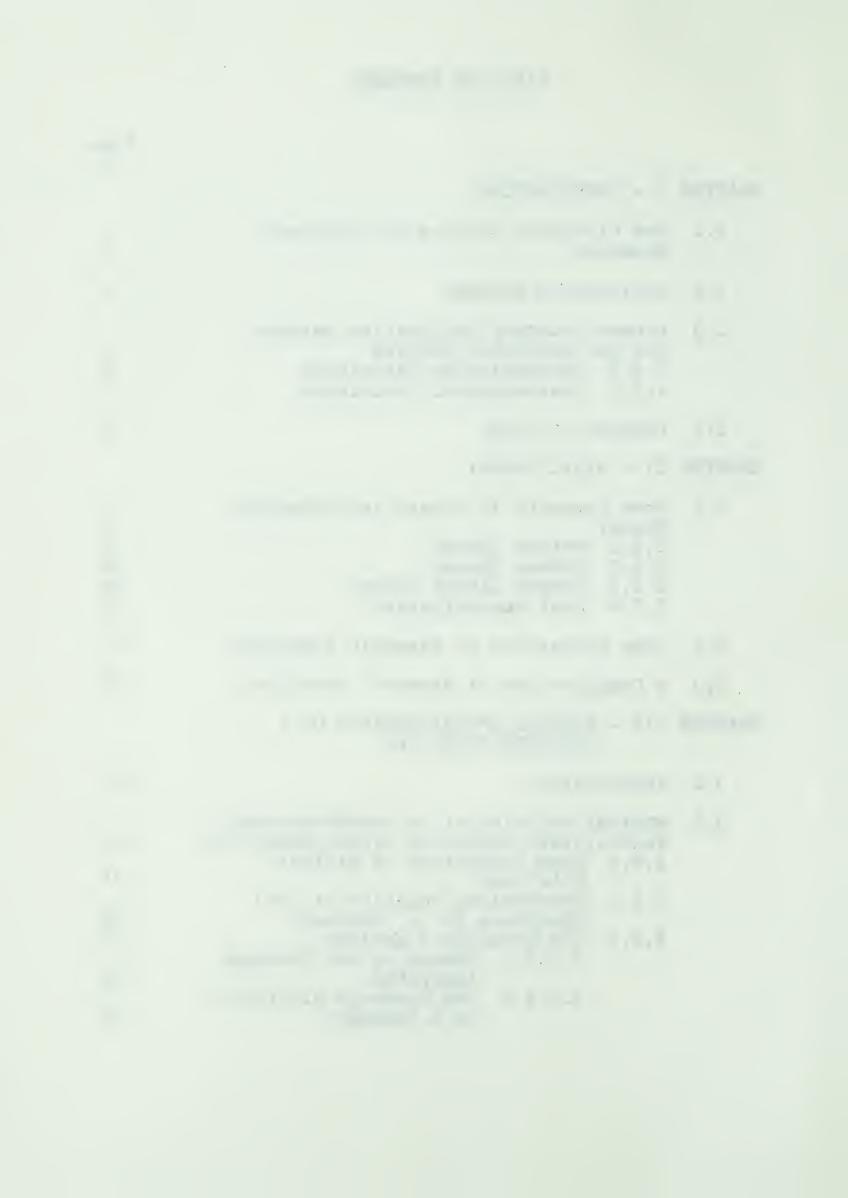reviewed, and numerical results are given.

ACKNOWLEDGEMENTS

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CHAPTER I

INTRODUCTION

This thesis reviews the exchange algorithm for computing a best linear approximation on a discrete point set and shows how it may be used by an algorithm of Remes to compute a best linear approximation on a continmum. These two algorithms are employed by a linear boundary collocation method to compute a close approximation to the solution of the Dirichlet problem for Laplace's equation.

1.1 The Dirichlet Problem for Laplace's Equation

Consider the following classical problem from the theory of partial differential equations: Let C be a simple closed curve, R the region interior to C , and f a continuous function on C . Find a function $u = u(x,y)$ having continuous second order derivatives in R satisfying

$$(1.1) \qquad u_{xx} + u_{yy} = 0 \quad \text{in } R ,$$

and

$$(1.2) \qquad u = f \quad \text{on } C .$$

The subscripts in (1.1) denote partial differentiation.

Equation (1.1) is known as Laplace's equation. Functions having continuous second order derivatives in R and satisfying (1.1) are called harmonic functions in R. Equation (1.2) is the Dirichlet boundary condition for Laplace's equation. Determining the solution u is the Dirichlet problem for Laplace's equation. Ahlfors ([1], pp. 175-199) has proven that u exists and is unique.

The Dirichlet problem arises in a wide variety of physical contexts, including heat transfer, hydrodynamics, electricity, and aerodynamics. Sample problems may be found in Forsythe and Wasow [16] and Fox [17].

Except for some simple boundaries such as the rectangle and the circle, where analytical solutions can be obtained in closed form, most boundaries require that some numerical method be used to obtain the solution. Most of the numerical work devoted to the Dirichlet problem has been concerned with replacing the derivatives in (1.1) by appropriate finite differences and then determining the solution at a finite set of points in the region R. However, this technique is difficult to apply to regions whose boundaries are irregular, and the resulting table of solution values is not convenient for use in subsequent numerical work. Attempts have therefore been made to derive methods which will not only yield the solution in closed form, but also be easy to apply to irregularly shaped boundaries.

1.2 Collocation Methods

Attempts to solve boundary-value problems in partial differential equations are often made using collocation methods. One assumes as an approximation to the solution $u(x,y)$ an expression of the form $g(x,y;c)$ which depends on a vector of parameters $c = (c_1,\ldots,c_n)$ such that for arbitrary values of c

a) the differential equation is already satisfied exactly, or

b) the boundary conditions are already satisfied exactly, or

c) g satisfies neither the differential equation nor the boundary conditions.

One then tries to determine c so that g satisfies

in method (a), the boundary conditions, or

in method (b), the differential equation, or

in method (c), the boundary conditions and the differential equation

as accurately as possible in some sense yet to be defined.

Definition 1.1. *A collocation method of type (a) described previously is called a boundary collocation method.*

For partial differential equations, boundary collocation methods are usually preferred for two reasons. First, it is generally easier to find functions which satisfy a particular differential equation than it is to find functions which

satisfy the boundary conditions for a particular boundary.
Second, boundary collocation methods are usually chosen
since their use, in so far as integration is concerned,
requires the evaluation of integrals around the boundary C
rather than throughout the region R . These observations
were also verified in practice by Shuleshko [29], who made
a comparative analysis of the three collocation techniques
applied to the torsion problem and concluded that not only
were boundary collocation methods easier to apply, but also
yielded better results.

Linear partial differential equations have the property
that $\alpha u + \beta v$ is a solution of the partial differential
equation if u and v satisfy the equation and α and β
are any two constants. Hence boundary collocation methods
for linear partial differential equations can have a
special form.

Definition 1.2. *A boundary collocation method
for a linear partial differential equation which uses an
approximating function of the form $g = \sum_{j=1}^{n} c_j g_j$, where
each function g_j , j=1,...,n , satisfies that partial
differential equation, is called a linear boundary
collocation method.*

1.3 Linear Boundary Collocation Methods for the Dirichlet Problem

Because Laplace's equation is linear, and since the sum of two continuous functions is continuous, the following is a linear boundary collocation method for the Dirichlet problem: Let $g_1,\ldots,g_n$ be harmonic functions in the region R. Then for an arbitrary vector of coefficients $c = (c_1,\ldots,c_n)$, $\sum_{j=1}^{n} c_j g_j$ is also a harmonic function in R. All linear boundary collocation methods are distinguished by the manner in which c is determined.

1.3.1 Interpolation Techniques

In this method, n distinct points $s_1,\ldots,s_n$ are selected on the boundary C, and the parameters are chosen so that $\sum_{j=1}^{n} c_j g_j$ interpolates to f at these points. That is, $(c_1,\ldots,c_n)$ is the solution of the system of linear equations,

$$(1.3) \qquad \sum_{j=1}^{n} c_j g_j(s_i) = f(s_i) , \quad i=1,\ldots,n .$$

The essential idea of this technique was first suggested by Walsh [34] in 1929, when he discussed the problem of approximating harmonic functions by harmonic polynomials. Since then the method has been applied with varying success to a number of engineering problems, including ones in heat transfer, elasticity, and the bending and buckling

of plates [4,5,6,7,8,19,23,26,28,30,33]. The technique
was first dubbed the "point-matching method" by Conway [6]
in 1961.

Not too much is known about this method theoretically.
Its use seems to indicate that its success depends upon the
geometry of the problem, but this dependence has not been
explicitly determined. Another problem arises from the
fact that the coefficient matrix in (1.3) may be singular
for some sets of points on C . Further, it is not known
under what conditions the method converges. Perhaps the
most critical problem is that a general set of properties
of the "best" set of matching points has not been found.

Theoretical material concerning this technique is
discussed in a series of papers by Curtiss [9-11] and
Sobczyk [31]. On the basis of simple cases of circles and
ellipses that have been investigated successfully [35,36],
and also arguing by analogy from the case of complex
analytic interpolation, Curtiss [9] concludes that for
sufficiently regular boundaries there probably exist
sequences of points for which this interpolation process
converges for a wide class of boundary data. However, his
suggestion that these points are probably obtainable as
the exterior conformal image of equally spaced points on
the unit circle (such as the roots of unity) is unsuitable
for numerical work because the conformal map is not always

available.   Sobczyk [31] has even shown that the resulting
image points on C may result in a singular coefficient
matrix in (1.3).

1.3.2 Least-Squares Techniques In a least-squares
approach, c is chosen to minimize

$$(1.4) \qquad E_1 = \int_C [f(s) - \sum_{j=1}^{n} c_j g_j(s)]^2 \, ds \; .$$

Usually the integration in (1.4) cannot be carried out
analytically. Therefore, m points are chosen on the
boundary C and m weights w_i corresponding to a
quadrature rule, and then the vector c is chosen to
minimize

$$(1.5) \qquad E_2 = \sum_{i=1}^{m} w_i[f(s_i) - \sum_{j=1}^{n} c_j g_j(s_i)]^2 \; .$$

In contrast to the interpolation technique, the discrete
process in (1.5) should be convergent as n becomes large
and m>>n . Details of this method and numerical results
can be found in papers by Lo [20], Davis [12], Davis and
Rabinowitz [14], O'Jalvo and Linzer [24], Merriman [21],
and Sparrow [32].

1.4 Purpose of Study

In Chapter II, we shall review some basic theory relevant to the Dirichlet problem and thus attempt to show that choosing c to minimize

$$(1.6) \qquad M = \max_{s \in C} \left| f(s) - \sum_{j=1}^{n} c_j g_j (s) \right|$$

is a more natural linear boundary collocation method than either the point-matching method or the least-squares method. In Chapter III, IV and V we shall show that such a vector c always exists and we shall present a computationally oriented algorithm for finding it. Results of numerical computations will be given in Chapter VI, and conclusions and suggestions for further research in Chapter VII.

A discrete version of this approach has been made by Quon and Leung [27] using linear programming techniques.

BASIC THEORY

2.1   Some Concepts in Linear Approximation Theory

The problem outlined in the previous chapter (see Section 1.4) is a typical one in approximation theory in that it involves the selection from a given set of functions one that is in some sense "close" to a prescribed function (usually) not in the set. Any discussion of approximation requires that we have a way of measuring the "closeness" of two functions. It is reasonable to demand that this measure should have some of the properties of geometrical distance, or more precisely, that it should be a metric.

2.1.1   Metric Spaces

Definition 2.1. Let S be a set of elements $x, y, z, \ldots,$ and let d be a real-valued function defined on pairs of elements of S . Then d is called a metric if it possesses the following properties:

$$
\begin{aligned}
&i) \quad d(x,x) = 0, \\
&ii) \quad d(x,y) > 0 \quad if \quad x \neq y, \\
&iii) \quad d(x,y) = d(y,x), \\
&iv) \quad d(x,y) \leq d(x,z) + d(z,y).
\end{aligned}
$$

(2.1)

The ordered pair (S,d) is then termed a metric space.

If, for example, S is the set of real numbers, the usual metric is $d(x,y) = |x-y|$. Properties (i) to (iv) are then simply familiar properties of the absolute value function.

The concept of a metric space is usually too general to yield many practical results and therefore the setting for many problems in the theory of approximation is a metric space of a particular kind called a normed linear space. We must first, however, define a linear space.

2.1.2 Linear Spaces

Definition 2.2. *Let* S *be a set of elements* $x,y,z,\ldots$ *for which two types of operations are possible.* S *is termed a linear space if any two elements* x,y *in* S *determine a unique element* $x+y$ *in* S *as their sum, if each element* x *in* S *and each scalar* α *determines a unique element* αx *in* S *as their scalar product, and if summation and scalar multiplication satisfy the following properties:*

    i) $x+y = y+x,$

    ii) $x + (y+z) = (x+y) + z,$

    iii)   *There exists a unique element* Θ *in* S
           such that $x+\Theta = x$ *for all* x *in* S
           (Θ *is called the zero element of the space),*

(2.2)   *iv) For each x in S , there exists a unique*

             inverse (-x) in S such that x + (-x) = Θ,

        v) α(βx) = (αβ)x for all scalars α,β,

      vi) α(x+y) = αx + αy,

     vii) (α+β)x = αx + βx,

    viii) 1x = x.

We now list some linear spaces which are important to this
study.

Example 2.1. *The real, n-dimensional, Euclidian*
space, denoted by R_n .
R_n consists of vectors x,y,z,... that are n-tuples of
real numbers: $x = (x_1,...,x_n)$. Summation is defined by
$x+y = (x_1+y_1,...,x_n+y_n)$, and scalar multiplication by
$αx = (αx_1,...,αx_n)$. The zero element is $Θ = (0,...,0)$,
and the inverse of x is $(-x) = (-x_1,...,-x_n)$.

Example 2.2. *The space of all functions*
continuous in the interval [a,b] , denoted by C[a,b] .
Let f,g,... be members of C[a,b] and let x be any
point in [a,b] . We define the sum of f and g by
$(f+g)(x) = f(x) + g(x)$, and the scalar product of α and
f by $(αf)(x) = αf(x)$. The zero element is that function
which vanishes identically in [a,b] and the inverse of
f is given by $(-f)(x) = -f(x)$.

Example 2.3. *The space of all functions continuous on a simple closed curve C , denoted by W(C) .*

Example 2.4. *Let C be a simple closed curve and R the region interior to C . The space of all complex functions analytic in R and continuous on C , denoted by A(R,C) .*

Before leaving the concept of a linear space, we shall present a definition which will be used in subsequent work.

Definition 2.3. *Let S be a linear space. A subset M of S is called convex if x,y are any two distinct elements of M implies that αx + (1-α)y is also a member of M for any real α , 0 ≤ α ≤ 1 .*

If we picture M geometrically, then αx + (1-α)y lies on the line segment joining x and y . Thus M is convex if it contains all the points on the line segment joining any two of its points. Circles, ellipses, rect-angles, and straight lines are examples of convex sets in R_n . The figure "8" is not convex in R_n .

2.1.3 Normed Linear Spaces

Definition 2.4. *A normed linear space is a linear space S for which there is defined a real-valued function on elements x,y,z,... in S called a norm, denoted by $\|x\|$, and satisfying the following laws:*

$$i) \quad \|\Theta\| = 0,$$

$$ii) \quad \|x\| > 0 \quad if \quad x \neq \Theta,$$

(2.3)

$$iii) \quad \|\alpha x\| = |\alpha| \, \|x\| \quad for \ any \ scalar \quad \alpha,$$

$$iv) \quad \|x+y\| \leq \|x\| + \|y\| \, .$$

An easily verified property of norms is given by

$$\underline{\text{Theorem 2.1.}} \quad \left| \, \|x\| - \|y\| \, \right| \leq \|x-y\| \, .$$

There are usually an infinity of norms that may be introduced into linear spaces in order to make them normed linear spaces. However, we discuss only those relevant to this study.

$\underline{\text{Theorem 2.2.}}$ *The linear space* R_n *with norm defined by*

(2.4)
$$\|x\| = \max_{1 \leq i \leq n} |x_i|$$

is a normed linear space.

Proof: We must verify that the function given in (2.4) satisfies the properties listed in (2.3). The first three properties are easily verified as simple consequences of properties of the absolute value function. To prove (iv) we write

$$\|x+y\| = \|(x_1+y_1,\ldots,x_n+y_n)\|$$

$$= \max_{1\le i\le n} |x_i+y_i|$$

$$= |x_k+y_k| \text{ , for some } k \text{ , } 1\le k\le n$$

$$\le |x_k| + |y_k|$$

$$\le \max_{1\le i\le n} |x_i| + \max_{1\le i\le n} |y_i|$$

$$= \|x\| + \|y\| \text{ .}$$

Theorem 2.3. *The linear space* $C[a,b]$ *with norm defined by*

$$(2.5) \qquad \|f\| = \max_{a\le x\le b} |f(x)|$$

is a normed linear space.

Theorem 2.4. *The linear space* $W(C)$ *with norm defined by*

$$(2.6) \qquad \|f\| = \max_{(x,y)\in C} |f(x,y)|$$

is a normed linear space.

Theorem 2.5. *The linear space $A(R,C)$ with norm*

defined by

(2.7) $$\|f\| = \max_{z \in C} |f(z)|$$

is a normed linear space.

Proof: The maximum modulus need only be taken over C, since a complex function analytic in R and continuous on C assumes its maximum modulus on C ([25], pp. 290-291).

Theorem 2.6. *In a normed linear space*

(2.8) $$d(f,g) = \|f-g\|$$

defines a metric.

Proof: We must show that (2.3) and (2.8) imply (2.1).

i) $d(f,f) = \|f-f\| = \|0 \cdot f\| = |0| \, \|f\| = 0$.

ii) $d(f,g) = \|f-g\| > 0$ since $f \neq g$ implies $f-g \neq \theta$.

iii) $d(f,g) = \|f-g\| = |-1| \, \|f-g\| = \|g-f\| = d(g,f)$

iv) $d(f,g) = \|f-g\| = \|f-h+h-g\| \leq \|f-h\| + \|h-g\|$

$$= d(f,h) + d(h,g) .$$

2.1.4 <u>Best Approximation</u> From the remarks made at the beginning of this Section and from Theorem 2.6, we see that a meaningful way of defining the closeness of two functions f and g in a normed linear space is to define it as the norm of their difference.

<u>Definition 2.5.</u> *Let* $f, g_1, \ldots, g_n$ *be elements of a normed linear space* S . *A best approximation to* f *by linear combinations of* $g_1, \ldots, g_n$ *is an element* $\sum\limits_{j=1}^{n} c_j g_j$ *for which*

$$(2.9) \qquad \left\| f - \sum_{j=1}^{n} c_j g_j \right\| \leq \left\| f - \sum_{j=1}^{n} a_j g_j \right\|$$

for every choice of constants $a_1, \ldots, a_n$.

Then the fundamental assertion of linear approximation theory is

<u>Theorem 2.7.</u> *The problem of finding* $\min\limits_{a_j} \left\| f - \sum\limits_{j=1}^{n} a_j g_j \right\|$ *has a solution.*

A proof of this important theorem can be found in ([13], pp. 137-139).

<u>Definition 2.6.</u> *For a given* $f, g_1, \ldots, g_n$ *in a normed linear space,* $E_n(f)$ *is defined by*

$$(2.10) \qquad E_n(f) = \min_{a_j} \left\| f - \sum_{j=1}^{n} a_j g_j \right\| .$$

Theorem 2.8.

(2.11) $$E_1(f) \geq E_2(f) \geq \ldots$$

Proof: (2.11) is true since linear combinations of $g_1, \ldots, g_n$ are also linear combinations of $g_1, \ldots, g_n, g_{n+1}$.

We have observed there is always one best approximation to f by linear combinations $\sum_{j=1}^{n} a_j g_j$. However, there may be more than one. If M is the set of best approximations, then we have

Theorem 2.9. *M is a convex set.*
A proof of the theorem can be found in ([13], pp. 140-141).

In our discussion of best approximation, the particular norm used was not specified. However, since the norms given in Theorems 2.2-2.5 are formulated in terms of a certain maximum, the corresponding best approximations solve the problem of minimizing the maximum error, and therefore we will subsequently refer to these best approximations as minimax approximations.

2.2 Some Properties of Harmonic Functions

Laplace's equation bears a close relationship to the theory of analytic functions. In fact, we have

Theorem 2.10. *The real part of any analytic function is a harmonic function; and conversely, every*

harmonic function is the real part of some analytic function.

Proof: Set $z = x + iy$ and let $f(z) = u(x,y) + iv(x,y)$ be an analytic function of z. Then u and v are related by the Cauchy-Riemann equations

$$(2.12) \qquad u_x = v_y \ , \ u_y = -v_x \ ,$$

and are said to be conjugate ([25], p. 83). Furthermore, u and v possess continuous partial derivatives of all orders ([25], pp. 185-186). If the first of these equations is differentiated with respect to x, the second with respect to y, and the two resulting equations added, then it can be seen that $u(x,y) = \text{Re } f(z)$ is a solution of Laplace's equation. On the other hand, let u be a harmonic function. Then equations (2.12) can be solved for v, and since the partial derivatives u_x, u_y, v_x, v_y are continuous, $f(z) = u(x,y) + iv(x,y)$ is an analytic function of z ([25], p. 85).

A parallel argument yields a similar theorem for the imaginary parts of analytic functions.

Another very important property of harmonic functions is given in

Theorem 2.11. *Let C be a simple closed curve and R the region interior to C. Let u be a harmonic*

function in R *and continuous on* C . *Then there exist*

two points z_1 *and* z_2 *on* C *such that* $u(z_1) \leq u(z) \leq u(z_2)$

for all z *in* R .

A proof of this theorem can be found in ([25], pp. 349-350).

A crucial result of this theorem is

Theorem 2.12. *Consider Dirichlet's problem for*

Laplace's equation (see Section 1.1). If $g(x,y)$ *is a*

harmonic function such that

$$(2.13) \qquad \max_{(x,y)\epsilon C} |f(x,y) - g(x,y)| \leq M$$

for some constant M , *then*

$$(2.14) \qquad \max_{(x,y)\epsilon R} |u(x,y) - g(x,y)| \leq M .$$

Proof: If u is the solution of the original

problem, then $u-g$ is the solution of the same problem with

boundary function $f-g$ rather than f . Note that $f-g$

is continuous and $u-g$ is harmonic. However, by Theorem 2.11,

a harmonic function assumes its maximum value on the boundary

C . That is,

$$\max_{(x,y)\epsilon R} |u(x,y) - g(x,y)| \leq \max_{(x,y)\epsilon C} |f(x,y) - g(x,y)| \leq M .$$

2.3 A Complete Set of Harmonic Functions

So far we have not specifically given a set of harmonic functions $g_1,...,g_n$ with which to approximate u in R and f on C. We see, from Theorem 2.10, that we need only use the real and/or imaginary parts of some analytic functions. Now z^n, $n=0,1,...$ are obviously analytic functions and hence

(2.15) $H = \{1,\ \text{Re}\ z^n,\ \text{Im}\ z^n,\ n=1,2,...\}$

is an infinite set of harmonic functions. They are commonly called harmonic polynomials.

Once a particular set of harmonic functions have been chosen, a problem arises as to the "completeness" of this set. Completeness refers to the possibility of arbitrarily close approximation of the boundary data by means of the boundary values of linear combinations of the selected set of harmonic functions. A relevant result, due to J.L. Walsh, is that the powers, $1,z,z^2,...$ are complete in $A(R,C)$ ([13], pp. 275-278). An important consequence of Walsh's assertion is

Theorem 2.13. *The harmonic polynomials are complete in* $W(C)$.

Proof: Let u be the solution of the Dirichlet problem. Then u is the real part of some analytic function $F(z)$, by Theorem 2.10. By the completeness of the powers $1, z, z^2, \ldots$ in $A(R,C)$, there exists a complex polynomial $p_n(z)$ such that $\max\limits_{z \in C} |F(z) - p_n(z)| \leq \delta$ for any $\delta > 0$. We now write

$$\delta \geq \max_{z \in C} |F(z) - p_n(z)|$$

$$\geq |F(z) - p_n(z)| \quad \text{for any } z \text{ on } C$$

$$\geq |\text{Re} (F(z) - p_n(z))|$$

$$\geq |f(x,y) - \text{Re } p_n(z)|$$

Since the above inequality is true for any z on C , we have that

$$\delta \geq \max_{(x,y) \in C} |f(x,y) - \text{Re } p_n(z)| .$$

Letting δ tend to zero proves the theorem.

CHAPTER III

MINIMAX APPROXIMATIONS ON A DISCRETE POINT SET

3.1 Introduction

Given the Dirichlet problem outlined in Chapter I (see Section 1.1), we set ourselves the task of selecting a vector of parameters $c = (c_1, \ldots, c_n)$ to minimize the expression

$$(3.1) \qquad M_1 = \max_{(x,y) \varepsilon C} \left| f(x,y) - \sum_{j=1}^{n} a_j g_j(x,y) \right| .$$

Our discussions in Chapter II ensured us that such a vector c existed. Furthermore, if g_j, $j=1,\ldots,n$ were chosen to be the harmonic polynomials, then for sufficiently large n, $\sum_{j=1}^{n} c_j g_j$ could be accepted as a reasonably good approximate solution to the Dirichlet problem.

As yet, however, we have not tackled the problem of how to determine the vector c. As a first step, we might select m points (x_i, y_i), $i=1,\ldots,m$ $(m>n)$ from the boundary C and minimize the expression

$$(3.2) \qquad M_2 = \max_{1 \leq i \leq m} \left| f(x_i, y_i) - \sum_{j=1}^{n} a_j g_j(x_i, y_i) \right| ,$$

accepting the minimax solution to (3.2) as a close approximation to the minimax solution of (3.1). This is equivalent to solving the following linear system of equations in the minimax sense:

$$(3.3) \qquad \sum_{j=1}^{n} a_j g_j(x_i, y_i) = f(x_i, y_i) , \quad i=1,\ldots,m .$$

Obtaining a minimax solution to an overdetermined, inconsistent system of linear equations is an important problem in its own right. The discussion we now give parallels that given by Handscomb ([18], pp. 73-82). A more advanced and complete approach has been given by Cheney ([3], pp. 28-56).

3.2 Minimax Solution of an Overdetermined, Inconsistent System of Linear Equations

Consider the system of linear equations

$$(3.4) \qquad \sum_{j=1}^{n} a_{ij} x_j = b_i , \quad i=1,\ldots,m .$$

We assume that the numbers a_{ij} and b_i are given and that the unknowns x_j are to be determined. (3.4) is frequently written in matrix notation as $Ax = b$. The system may have exactly one solution, infinitely many

solutions, or no solution, depending on the data. If the
rank of the coefficient matrix A is equal to the rank
of the augmented matrix A_b , then (3.4) possesses a
solution and is said to be consistent. If m>n and the
rank of A does not equal the rank of A_b , then (3.4)
is said to be overdetermined and inconsistent. However,
even though (3.4) may possess no solution, we may still
try to find a vector $x = (x_1, \ldots, x_n)$ that minimizes the
expression

(3.5)
$$\rho = \max_{1 \leq i \leq m} \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| .$$

Such an x is called a minimax solution of (3.4).

3.2.1 Some Properties of Minimax Solutions

Theorem 3.1. *Let values* a_{ij}, b_i *be given for*
j=1,...,n and i=1,...,m>n . The problem of finding

$$\min_{x \in R_n} \max_{1 \leq i \leq m} \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right|$$

has a solution.

Proof: For each j=1,...,n , let A_j denote the
column vector $(a_{1j}, \ldots, a_{mj})^T$. Then (3.4) can be rewritten

in the form

$$(3.6) \qquad \sum_{j=1}^{n} x_j A_j = b \; .$$

Now $\sum_{j=1}^{n} y_j A_j$ and b are elements of R_m for any

$y = (y_1, \ldots, y_n)$. Hence, by Theorems 2.2 and 2.7, there

exists a vector $x \varepsilon R_n$ such that

$$(3.7) \qquad \| b - \sum_{j=1}^{n} x_j A_j \| \leq \| b - \sum_{j=1}^{n} y_j A_j \|$$

for any $y \varepsilon R_n$.

Theorem 3.2. *The set of minimax xolutions to
(3.4) is convex.*

Proof: Theorem 3.2 is just a rewording of
Theorem 2.9.

Before proceeding further, we make the following
important definition.

Definition 3.1. *A set of vectors in R_n is said
to satisfy Haar's condition if every subset of n of them
is linearly independent.*

We now make the rather strong assumption that the set
of row vectors making up the matrix A satisfy Haar's
condition. Therefore, any n of the m equations in (3.4)

possess a nonsingular coefficient matrix. Haar's condition
will be discussed in detail in Chapter IV.

   <u>Theorem 3.3</u>. *Let* $x = (x_1, \ldots, x_n)$ *be a minimax
solution to (3.4) and let* $\rho = \max\limits_{1 \le i \le m} |b_i - \sum\limits_{j=1}^{n} a_{ij}x_j|$.
Then $I = \{i : |b_i - \sum\limits_{j=1}^{n} a_{ij}x_j| = \rho\}$ *contains at least
(n+1) elements.*

   <u>Proof</u>: Suppose, on the contrary, that I contains
r elements, where $r < n+1$. Reordering the equations if
necessary, we can assume that the maximum absolute error
occurs in the first r equations. Thus we have

$$b_1 - \sum_{j=1}^{n} a_{1j}x_j = s_1\rho$$

$$\cdots \cdots \cdots \cdots$$

$$b_r - \sum_{j=1}^{n} a_{rj}x_j = s_r\rho$$

(3.8)

$$b_{r+1} - \sum_{j=1}^{n} a_{r+1,j}x_j = s_{r+1}e_{r+1}$$

$$\cdots \cdots \cdots \cdots$$

$$b_m - \sum_{j=1}^{n} a_{mj}x_j = s_m e_m$$

where

(3.9) $\qquad 0 \leq e_i < \rho$, $i = r+1, \ldots, m$,

and

(3.10) $\qquad s_i = \text{sgn}(b_i - \sum_{j=1}^{n} a_{ij} x_j)$, $i = 1, \ldots, m$.

We now show that the rank of the matrix $A_r = (a_{ij})$, $j=1, \ldots, n$; $i=1, \ldots, r$; is less than r . If, on the other hand, the rank of A_r equals r , then we can find a vector $u = (u_1, \ldots, u_n)$ such that

(3.11) $\qquad \sum_{j=1}^{n} a_{ij} u_j = s_i \rho$, $i = 1, \ldots, r$.

Now define a new vector $y = (y_1, \ldots, y_n)$ by

(3.12) $\qquad y_j = x_j + \lambda u_j$, $j = 1, \ldots, n$,

where

(3.13) $\qquad 0 < \lambda < \min(\dfrac{\rho - M_1}{M_2} , 1)$.

M_1 and M_2 in (3.13) are given by

$$(3.14) \qquad M_1 = \max_{r+1 \leq i \leq m} e_i \, ,$$

and

$$(3.15) \qquad M_2 = \max_{r+1 \leq i \leq m} \left| \sum_{j=1}^{n} a_{ij} u_j \right| \, .$$

Then for $i=1,\ldots,m$ we have

$$b_i - \sum_{j=1}^{n} a_{ij} y_j = b_i - \sum_{j=1}^{n} a_{ij}(x_j + \lambda u_j)$$

$$= b_i - \sum_{j=1}^{n} a_{ij} x_j - \lambda \sum_{j=1}^{n} a_{ij} u_j$$

$$(3.16)$$

$$= \begin{cases} s_i \rho - \lambda \, s_i \rho \, , & i=1,\ldots,r \, . \\[2em] s_i \rho_i - \lambda \sum_{j=1}^{n} a_{ij} u_j \, , & i=r+1,\ldots,m \, . \end{cases}$$

Taking the absolute value of (3.16) we get

(3.17)
$$|b_i - \sum_{j=1}^{n} a_{ij}y_j| = \begin{cases} (1-\lambda)\rho \ , & i=1,\ldots,r \ . \\ \\ |s_ie_i - \lambda \sum_{j=1}^{n} a_{ij}u_j| \ , & i=r+1,\ldots,m \ . \end{cases}$$

Certainly, $(1-\lambda)\rho < \rho$ and for $i=r+1,\ldots,m$

$$|s_ie_i - \lambda \sum_{j=1}^{n} a_{ij}u_j| \leq e_i + \lambda \, |\sum_{j=1}^{n} a_{ij}u_j|$$

$$< e_i + \frac{\rho-M_1}{M_2} \, |\sum_{j=1}^{n} a_{ij}u_j|$$

(3.18)

$$\leq e_i + \rho - M_1 \ , \text{ by (3.15)}$$

$$\leq \rho \ , \text{ by (3.14)}.$$

This means that y is a better minimax approximation than x , which is a contradiction. Hence, the rank of A_r is less than r . Since $r \leq n$, this contradicts our basic assumption that the rows of A satisfy Haar's condition. Hence, I contains at least $(n+1)$ elements.

Theorem 3.4. *The minimax solution to (3.4) is unique.*

Proof: Suppose that $x, y \in R_n$ are two distinct minimax solutions and let

$$(3.19) \quad \rho = \max_{1 \leq i \leq m} \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| = \max_{1 \leq i \leq m} \left| b_i - \sum_{j=1}^{n} a_{ij} y_j \right| .$$

Now define

$$(3.20) \quad v = \tfrac{1}{2}(x+y) .$$

Then v is also a minimax solution by Theorem 3.2. For those i for which

$$(3.21) \quad b_i - \sum_{j=1}^{n} a_{ij} v_j = s_i \rho ,$$

where s_i is still given by (3.10), we must have

$$(3.22) \quad \left(b_i - \sum_{j=1}^{n} a_{ij} x_j \right) + \left(b_i - \sum_{j=1}^{n} a_{ij} y_j \right) = 2 s_i \rho .$$

Therefore, we obtain

$$(3.23) \qquad 2\rho \leq \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| + \left| b_i - \sum_{j=1}^{n} a_{ij} y_j \right| .$$

Neither term on the right of the above inequality can be greater than ρ. Therefore,

$$(3.24) \qquad 2\rho = \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| + \left| b_i - \sum_{j=1}^{n} a_{ij} y_j \right| .$$

By the same reasoning, we obtain

$$(3.25) \qquad \rho = \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| = \left| b_i - \sum_{j=1}^{n} a_{ij} y_j \right| .$$

Hence,

$$(3.26) \qquad b_i - \sum_{j=1}^{n} a_{ij} x_j = b_i - \sum_{j=1}^{n} a_{ij} y_j ,$$

for if the two terms were of opposite sign, this would contradict (3.22). Therefore,

$$(3.27) \qquad \sum_{j=1}^{n} a_{ij} (x_j - y_j) = 0 .$$

By Theorem 3.3, there are $(n+1)$ such i for which (3.27) holds. Choose any n of them. The above equation is then a system of n equations in n unknowns. Our basic assumption ensures that the coefficient matrix of this system is nonsingular. Hence, (3.27) has only the zero solution. That is,

$$(3.28) \qquad x_j - y_j = 0 \ , \quad j=1,\ldots,n \ ,$$

and therefore $x=y$ as was to be shown.

Now that we have some insight into the existence, characterization, and uniqueness of the minimax solution to (3.4), we must approach the problem of actually computing it. We first turn our attention to the special case where the number of equations is one greater than the number of unknowns.

3.2.2 The Minimax Solution of $(n+1)$ Equations in n Unknowns

From Theorem 3.3, we know the minimax solution yields the same absolute error in each equation. Call this absolute error ρ . Using the notation $s_i = \pm 1$ for the sign of the error in the ith equation, we have that the minimax solution x satisfies

$$(3.29) \qquad b_i - \sum_{j=1}^{n} a_{ij} x_j = s_i \rho \ , \quad i=1,\ldots,n+1 \ .$$

For a particular choice of signs, the above system can be rewritten as

$$(3.30) \qquad \sum_{j=1}^{n} a_{ij} x_j + s_i \rho = b_i \ , \quad i=1,\ldots,n+1 \ ,$$

which is simply a system of $(n+1)$ equations in the $(n+1)$ unknowns $x_1,\ldots,x_n,\rho$. Let us define

$$(3.31) \quad T_i = (-1)^i \begin{vmatrix} a_{11} & \cdots & a_{1n} \\ \vdots & & \vdots \\ a_{i-1,1} & \cdots & a_{i-1,n} \\ a_{i+1,1} & \cdots & a_{i+1,n} \\ \vdots & & \vdots \\ a_{n+1,1} & \cdots & a_{n+1,n} \end{vmatrix} \ , \quad i=1,\ldots,n+1 \ .$$

Note that by Haar's condition, $T_i \neq 0$ for all i . We then have

Theorem 3.5. *The signs of* $s_1,\ldots,s_{n+1}$ *are the same as or opposite to the corresponding signs of* $T_1,\ldots,T_{n+1}$.

Proof: The coefficient matrix of (3.30) is given by

$$
\begin{bmatrix}
a_{11} & \cdots & a_{1n} & s_1 \\
\vdots & & \vdots & \vdots \\
a_{n+1,1} & \cdots & a_{n+1,n} & s_{n+1}
\end{bmatrix}
$$

Solving for ρ by Cramer's rule, expanding determinants down the last column, and utilizing (3.31), we get

$$
(3.32) \qquad \rho = \left(\sum_{i=1}^{n+1} T_i b_i \right) \Big/ \sum_{i=1}^{n+1} s_i T_i .
$$

We are trying to determine s_i so that ρ is minimized. Hence, we must try to make the denonimator in (3.32) as large in absolute value as possible. This will happen if

$$
(3.33) \qquad s_i = \text{sgn } T_i , \quad i=1,\ldots,n+1 ,
$$

or

$$
(3.34) \qquad s_i = - \text{sgn } T_i , \quad i=1,\ldots,n+1 ,
$$

which was to be shown.

Now that the signs s_i have been independently determined, it remains to be shown that the coefficient matrix of (3.30) is nonsingular.

Theorem 3.6

(3.35)

$$D = \begin{vmatrix} a_{11} & \cdots & a_{1n} & s_1 \\ \vdots & & \vdots & \vdots \\ a_{n+1,1} & \cdots & a_{n+1,n} & s_{n+1} \end{vmatrix} \neq 0 .$$

Proof: We have already remarked that each $T_i \neq 0$. To prove the assertion, we expand the determinant down the last column, obtaining

(3.36)

$$D = \sum_{i=1}^{n+1} s_i T_i .$$

But by Theorem 3.5, we have chosen s_i according to (3.33) or (3.34), and in one case, we have a sum of all positive terms, and in the other, a sum of all negative terms. In either case,

(3.37)

$$D \neq 0 ,$$

and the theorem is established.

To summarize, we can obtain the minimax solution of (n+1) equations in n unknowns by first determining the numbers s_i according to (3.33) or (3.34) and then solving the linear system (3.30) for $x_1, \ldots, x_n, \rho$.

We now return to the original problem of determining the minimax solution to (3.4). The results just presented are important because the minimax solution of (3.4) is also the minimax solution to some subsystem of (3.4) consisting of just (n+1) equations (see Theorem 3.3). The main problem now, then, is to determine which subset of (n+1) equations is pertinent. This subset can be determined by the exchange algorithm.

3.2.3 <u>The Exchange Algorithm</u> The basic idea of this algorithm is to calculate the minimax solutions of a succession of subsystems, each consisting of (n+1) equations from (3.4). By Theorem 3.3, the minimax solution of one of these subsystems is the required minimax solution. Since there are only a finite number of possible sub-systems, the exchange algorithm will converge in a finite number of steps, providing we can show that it never reconsiders a previous subsystem.

3.2.3.1 <u>Theory of the Exchange Algorithm</u>
Initially, we choose any (n+1) equations from (3.4). Without loss of generality, we may assume that we have

selected the first $(n+1)$. Then form the T_i and s_i according to

$$(3.38) \quad T_i = (-1)^i \begin{vmatrix} a_{11} & \cdots & a_{1n} \\ \vdots & & \vdots \\ a_{i-1,1} & \cdots & a_{i-1,n} \\ a_{i+1,1} & \cdots & a_{i+1,n} \\ \vdots & & \vdots \\ a_{n+1,1} & \cdots & a_{n+1,n} \end{vmatrix} , \quad i=1,\ldots,n+1 \ ,$$

and

$$(3.39) \qquad s_i = \text{sgn } T_i \ , \quad i=1,\ldots,n+1 \ .$$

The values of the $(n+1)$ unknowns $x_1,\ldots,x_n,\rho$ are obtained by solving

$$(3.40) \qquad \sum_{j=1}^{n} a_{ij} x_j + s_i \rho = b_i \ , \quad i=1,\ldots,n+1 \ .$$

We may at this point assume that $\rho > 0$. If (3.40) gives $\rho < 0$, we can simply replace each s_i by $-s_i$ and each T_i by $-T_i$. We now compute the discrepancies in each

equation from

$$(3.41) \qquad \rho_i = \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| , \quad i=1,\ldots,m ,$$

and determine $\max_{1 \le i \le m} \rho_i$. A test is made to see if $\max_{1 \le i \le m} \rho_i > \rho$. If not, then x is our required minimax solution and ρ is the associated minimax error. On the other hand, if the above condition is true, then we may assume that $\max_{1 \le i \le m} \rho_i = \rho_{n+2}$, so that

$$(3.42) \qquad \sum_{j=1}^{n} a_{n+2,j} x_j + s_{n+2} \rho_{n+2} = b_{n+2} , \quad \rho_{n+2} > \rho ,$$

where, as usual, s_{n+2} is the sign of the error in the (n+2)nd equation.

Theorem 3.7. *If (3.42) holds, then the minimax error ρ' of some subsystem consisting of $(n+1)$ equations selected from equations $1,\ldots,n+2$ is greater than ρ .*

The method of proof will be to explicitly give the indices of the equations comprising the relevant subsystem.

Proof: We define a set of row vectors by

$$(3.43) \qquad A_i = (a_{i1},\ldots,a_{in}) , \quad i=1,\ldots,n+2 .$$

Now consider the determinants given by

$$(3.44) \qquad D_j = \begin{vmatrix} a_{1j} & a_{11} & \cdots & a_{1n} \\ \vdots & \vdots & & \vdots \\ a_{n+1,j} & a_{n+1,1} & \cdots & a_{n+1,n} \end{vmatrix} , \quad j=1,\ldots,n .$$

Obviously, $D_j = 0$ for all j since the 1st and $(j+1)^{st}$ columns of the determinant are identical. If we expand D_j down the first column, we obtain

$$(3.45) \qquad 0 = \sum_{i=1}^{n+1} a_{ij} T_i , \quad j=1,\ldots,n ,$$

which may be rewritten as

$$(3.46) \qquad \sum_{i=1}^{n+1} T_i A_i = 0 .$$

Furthermore, since $s_i^2 = 1$, we may rewrite (3.46) as

$$(3.47) \qquad \sum_{i=1}^{n+1} T_i s_i (s_i A_i) = 0 .$$

If we consider (3.47) as a linear combination of the vectors $s_i A_i$, then, by (3.39), all the coefficients are positive.

Since no $T_i = 0$, we may solve for any A_k from (3.46), getting

$$(3.48) \qquad A_k = \sum_{\substack{i=1 \\ i \neq k}}^{n+1} - \frac{T_i}{T_k} A_i \; .$$

Now by Haar's condition, we know that the vectors $A_1, \ldots, A_{n+1}$ span R_n. Therefore, there exist numbers λ_i such that

$$(3.49) \qquad A_{n+2} = \sum_{i=1}^{n+1} \lambda_i A_i \; .$$

Hence,

$$0 = A_{n+2} - \sum_{i=1}^{n+1} \lambda_i A_i$$

$$= A_{n+2} - \lambda_k A_k - \sum_{\substack{i=1 \\ i \neq k}}^{n+1} \lambda_k A_k$$

$$(3.50) \qquad = A_{n+2} - \lambda_k \sum_{\substack{i=1 \\ i \neq k}}^{n+1} - \frac{T_i}{T_k} A_i - \sum_{\substack{i=1 \\ i \neq k}}^{n+1} \lambda_i A_i$$

$$= A_{n+2} + \sum_{\substack{i=1 \\ i \neq k}}^{n+1} (\lambda_k \frac{T_i}{T_k} - \lambda_i) A_i$$

$$= s_{n+2} A_{n+2} + \sum_{\substack{i=1 \\ i \neq k}}^{n+1} s_{n+2} s_i (\lambda_k \frac{T_i}{T_k} - \lambda_i) s_i A_i \; .$$

The last line comes from multiplying the previous one by s_{n+2} and again noting that $s_i^2 = 1$. The above equation expresses the fact that we have found numbers such that a linear combination of $s_i A_i$, $i=1,\ldots,n+2$; $i \neq k$ with these numbers as coefficients is equal to zero. Referring to the discussion after (3.47), we must have all of these coefficients positive. That is, we require

$$(3.51) \qquad s_{n+2} s_i \left(\lambda_k \frac{T_i}{T_k} - \lambda_i \right) > 0 , \quad i=1,\ldots,n+1; \; i \neq k .$$

Therefore, we must have

$$(3.52) \qquad s_{n+2} s_i T_i \frac{\lambda_k}{T_k} > s_{n+2} s_i \lambda_i , \quad i=1,\ldots,n+1; \; i \neq k .$$

From (3.39), we have $s_i T_i > 0$. Hence,

$$(3.53) \qquad s_{n+2} \frac{\lambda_k}{T_k} > s_{n+2} \frac{\lambda_i}{T_i} , \quad i=1,\ldots,n+1; \; i \neq k ,$$

and therefore k is the index satisfying

$$(3.54) \qquad s_{n+2} \frac{\lambda_k}{T_k} = \max_{1 \leq i \leq n+1} s_{n+2} \frac{\lambda_i}{T_i} .$$

The index k thus chosen is unique, for otherwise one of
the coefficients in (3.50) would vanish, contradicting
Haar's condition. It is now asserted that the minimax
error ρ' associated with equations $1,\ldots,k-1,k+1,\ldots,n+2$
is greater than the minimax error ρ associated with
equations $1,\ldots,n+1$. Let x' be the minimax solution
of equations $1,\ldots,k-1,k+1,\ldots,n+2$. Then

$$(3.55) \qquad \left| b_i - \sum_{j=1}^{n} a_{ij} x_j \right| < \left| b_{n+2} - \sum_{j=1}^{n} a_{n+2,j} x_j \right| ,$$

$$i = 1, \ldots, n+1 ,$$

and

$$(3.56) \qquad \left| b_i - \sum_{j=1}^{n} a_{ij} x_j' \right| = \left| b_{n+2} - \sum_{j=1}^{n} a_{n+2,j} x_j' \right| ,$$

$$i = 1, \ldots, n+1; \ i \neq k .$$

Hence, $x \neq x'$, or

$$(3.57) \qquad\qquad x - x' \neq 0 .$$

Now we can write

$$s_i \sum_{j=1}^{n} a_{ij}(x_j'-x_j) = s_1(b_1 - \sum_{j=1}^{n} a_{1j}x_j)$$

(3.58)
$$- s_1(b_1 - \sum_{j=1}^{n} a_{1j}x_j')$$

$$= \rho - \rho' , \quad i=1,\ldots,n+1; \ i \neq k .$$

Now by Haar's condition and (3.57), we have

(3.59)
$$\rho - \rho' \neq 0 .$$

We also have that

$$s_{n+2} \sum_{j=1}^{n} a_{n+2,j}(x_j'-x_j) = s_{n+2}(b_{n+2} - \sum_{j=1}^{n} a_{n+2,j}x_j)$$

(3.60)
$$- s_{n+2}(b_{n+2} - \sum_{j=1}^{n} a_{n+2,j}x_j')$$

$$> \rho - \rho' .$$

We are trying to show that $\rho - \rho' < 0$. We have already shown that $\rho - \rho' \neq 0$. Suppose, then, that $\rho - \rho' > 0$. Then by (3.58) and (3.60), we can write

$$(3.61) \qquad s_i \sum_{j=1}^{n} a_{ij}(x_j'-x_j) > 0 , \quad i=1,\ldots,n+2; \; i \neq k .$$

From (3.50) and (3.54), we know there exist numbers μ_i such that

$$(3.62) \qquad \sum_{\substack{i=1 \\ i \neq k}}^{n+2} \mu_i s_i A_i = 0 ,$$

or equivalently,

$$(3.63) \qquad \sum_{\substack{i=1 \\ i \neq k}}^{n+2} \mu_i s_i a_{ij} = 0 , \quad j=1,\ldots,n .$$

Since the $\mu_i > 0$, we can rewrite (3.61) as

$$(3.64) \qquad \mu_i s_i \sum_{j=1}^{n} a_{ij}(x_j'-x_j) > 0 , \quad i=1,\ldots,n+2; \; i \neq k .$$

For any numbers y_j , (3.63) can be altered to read

$$(3.65) \qquad y_j \sum_{\substack{i=1 \\ i \neq k}}^{n+2} \mu_i s_i A_i = 0 , \quad j=1,\ldots,n .$$

Summing the last set of equations over j , we obtain

$$(3.66) \qquad \sum_{\substack{i=1 \\ i \neq k}}^{n+2} \mu_i s_i \sum_{j=1}^{n} y_j a_{ij} = 0 .$$

The above equation is true for any set of numbers y_j . In particular, then, it must be true for $y_j = x'_j - x_j$, $j=1,\ldots,n$, and thus we get

$$(3.67) \qquad \sum_{\substack{i=1 \\ i \neq k}}^{n+2} \mu_i s_i \sum_{j=1}^{n} a_{ij}(x'_j - x_j) = 0 .$$

If we sum (3.64) over all permissable indices of i , we get

$$(3.68) \qquad \sum_{\substack{i=1 \\ i \neq k}}^{n+2} \mu_i s_i \sum_{j=1}^{n} a_{ij}(x'_j - x_j) > 0 ,$$

which contradicts (3.67). Therefore, $\rho - \rho' < 0$ and thus $\rho' > \rho$ as was to be shown.

We can now state that the exchange algorithm converges in a finite number of steps. From the previous theorem, we know that the minimax error strictly increases during each cycle of the algorithm. Hence, the algorithm never returns to a previously considered subsystem. Since there are only a finite number of such subsystems, the algorithm must eventually reach the critical one whose existence is guaranteed by Theorem 3.3.

3.2.3.2 The Exchange Algorithm on a Computer

Our approach in this section will be to present the exchange algorithm step by step with explanations inserted where necessary.

From the last section, we know that we will always be considering some subset of (n+1) equations. This can be most easily handled by an (n+1) - component vector $I = \{i_1, i_2, \ldots, i_{n+1}\}$ which gives the indices of the (n+1) equations presently under consideration. We suppose that the matrix A , the vector b , and a vector I have been given. The exchange algorithm then proceeds as follows:

1)

$$T_k \leftarrow (-1)^k \begin{vmatrix} a_{i_1,1} & \cdots & a_{i_1,n} \\ \vdots & & \vdots \\ a_{i_{k-1},1} & \cdots & a_{i_{k-1},n} \\ a_{i_{k+1},1} & \cdots & a_{i_{k+1},n} \\ \vdots & & \vdots \\ a_{i_{n+1},1} & \cdots & a_{i_{n+1},n} \end{vmatrix} , \quad k=1,\ldots,n+1 .$$

Calculating the determinant is a standard problem in linear algebra. For this and all other problems in linear algebra we have adapted a set of programs given by Forsythe and Moler ([15], pp. 68-72). Their main computational tool is Gaussian elimination, with partial pivoting and iterative improvement. For completeness, these programs have been listed in the Appendix.

2) $s_k \leftarrow \mathrm{sgn}\, T_k , \quad k=1,\ldots,n+1 .$

3)
$$\begin{bmatrix} d_{11} & \cdots & d_{1,n+1} \\ \vdots & & \vdots \\ d_{n+1,1} & \cdots & d_{n+1,n+1} \end{bmatrix} \leftarrow \begin{bmatrix} a_{i_1,1} & \cdots & a_{i_1,n} & s_1 \\ \vdots & & \vdots & \vdots \\ a_{i_{n+1},1} & \cdots & a_{i_{n+1},n} & s_{n+1} \end{bmatrix}^{-1}$$

4)
$$\begin{bmatrix} x_1 \\ \vdots \\ x_n \\ \rho \end{bmatrix} \leftarrow \begin{bmatrix} d_{11} & \cdots & d_{1,n+1} \\ \vdots & & \vdots \\ d_{n+1,1} & \cdots & d_{n+1,n+1} \end{bmatrix} \begin{bmatrix} b_{i_1} \\ \vdots \\ b_{i_{n+1}} \end{bmatrix}$$

The accumulation of inner products is carried out in double precision.

5) $\quad \rho_i \leftarrow b_i - \sum\limits_{j=1}^{n} a_{ij} x_j$, $\quad i = 1, \ldots, m$.

6) $\quad |\rho_\alpha| = \max\limits_{1 \leq i \leq m} |\rho_i|$.

7) We test: Is $|\rho_\alpha| > \rho$? If the answer is no, then $[x_1, \ldots, x_n, \rho]$ gives the minimax solution and the associated minimax error. If the answer is yes, we have found it necessary to make a further test to see if $\alpha \, \epsilon \, I$. If so, then we exit with the vector $[x_1, \ldots, x_n, \rho]$ as before. If not, then we proceed to the next step.

8) $\quad u \leftarrow \text{sgn } \rho_\alpha$.

A word of explanation is necessary to explain the following steps in the algorithm. From the proof of Theorem 3.7, we see that we must express $s_\alpha A_\alpha$ as a linear combination of $s_1 A_{i_1}, \ldots, s_{n+1} A_{i_{n+1}}$. To do this, set

$$(3.69) \qquad A_\alpha = \sum_{k=1}^{n+1} \lambda_k A_{i_k} .$$

The above equation can be written in matrix form

$$(3.70) \qquad [\lambda_1, \ldots, \lambda_{n+1}] \begin{bmatrix} a_{i_1,1} & \cdots & a_{i_1,n} & s_1 \\ \vdots & & \vdots & \vdots \\ a_{i_{n+1},1} & \cdots & a_{i_{n+1},n} & s_{n+1} \end{bmatrix}$$

$$= [a_{\alpha 1}, \ldots, a_{\alpha n}, u] .$$

Therefore, the next step in the algorithm is given by

$$9) \qquad [\lambda_1, \ldots, \lambda_{n+1}] = [a_{\alpha 1}, \ldots, a_{\alpha n}, u] \begin{bmatrix} d_{11} & \cdots & d_{1,n+1} \\ \vdots & & \vdots \\ d_{n+1,1} & \cdots & d_{n+1,n+1} \end{bmatrix}$$

We must now determine k from (3.54). From the definition of an inverse, we have

$$\begin{bmatrix} d_{11} & \cdots & d_{1,n+1} \\ \vdots & & \vdots \\ d_{n+1,1} & \cdots & d_{n+1,n+1} \end{bmatrix} \begin{bmatrix} a_{i_1,1} & \cdots & a_{i_1,n} & s_1 \\ \vdots & & \vdots & \vdots \\ a_{i_{n+1},1} & \cdots & a_{i_{n+1},n} & s_{n+1} \end{bmatrix} =$$

(3.71)
$$\begin{bmatrix} 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \vdots & & \vdots \\ 0 & 0 & \cdots & 1 \end{bmatrix}$$

and hence,

(3.72)
$$\sum_{k=1}^{n+1} d_{n+1,k} A_{i_k} = 0 \ .$$

Comparing (3.72) with (3.46), we see that after the initial inverse has been calculated, the numbers $d_{n+1,k}$ may serve as the numbers T_k . From (3.69) we have

$$(3.73) \qquad u \, A_\alpha = \sum_{k=1}^{n+1} \lambda_k u s_k (s_k A_{i_k})$$

and therefore, the next step in the algorithm can be written

$$10) \quad u \, \frac{\lambda_k}{d_{n+1,k}} = \max_{1 \le i \le n+1} u \, \frac{\lambda_i}{d_{n+1,i}} \, .$$

The only thing left to do is to update the inverse. We are changing the kth row of a coefficient matrix by replacing it by a known linear combination of rows, the coefficients being the λ_k of (3.69). The effect on the inverse is given by Handscomb ([18], pp. 79-80).

$$11) \quad d_{ij} \leftarrow \begin{cases} d_{ij}/\lambda_j \, , & j=k; \ i=1,\ldots,n+1 \, . \\[2em] d_{ij} - \dfrac{\lambda_j}{\lambda_k} d_{ik} \, , & j \ne k; \ i=1,\ldots,n+1 \, . \end{cases}$$

$$12) \quad i_k \leftarrow \alpha \, .$$

The algorithm now returns to step 4.

A program to compute the minimax solution of (3.4) is given in the Appendix.

CHAPTER IV

HAAR'S CONDITION

4.1 An Example

In the previous chapter we presented a detailed exposition of the exchange algorithm for computing the minimax solution of an overdetermined, inconsistent system of linear equations $Ax = b$ (see Section 3.4). The theoretical basis of this algorithm relied heavily on the assumption that the row vectors comprising the coefficient matrix A satisfied Haar's condition (see Definition 3.1). To see the difficulties arising when Haar's condition is violated, consider the problem of determining the minimax solution of the following system of linear equations

$$2x_1 + x_2 = 6.9$$

$$3x_1 + x_2 = 7.2$$

$$x_1 + x_2 = 3.0$$

(4.1)

$$x_1 - x_2 = 1.0$$

$$2x_1 + 4x_2 = 11.1$$

$$x_1 + 2x_2 = 7.0$$

We note that Haar's condition is violated since the vectors
(2,4) and (1,2) corresponding to the fifth and sixth
equations are linearly dependent. At each stage, the
exchange algorithm considers subsystems of three equations.
As in Section 3.2.3.2, we will let I be a 3-component
vector giving the indices of the equations presently under
consideration. Initially, suppose I = [1,2,3] . We
present the computations of the exchange algorithm given
in Section 3.2.3.2.

Step 2

1) $T_1 = (-1) \begin{vmatrix} 3 & 1 \\ 1 & 1 \end{vmatrix} = -2, T_2 = 1, T_3 = 1$.

2) $s_1 = \text{sgn } T_1 = -1, s_2 = 1, s_3 = 1$.

3)
$$D = \begin{bmatrix} 2 & 1 & -1 \\ 3 & 1 & 1 \\ 1 & 1 & 1 \end{bmatrix}^{-1} = \begin{bmatrix} 0 & \frac{1}{2} & -\frac{1}{2} \\ \frac{1}{2} & -\frac{3}{4} & \frac{5}{4} \\ -\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \end{bmatrix}$$

4) $$\begin{bmatrix} x_1 \\ x_2 \\ \rho \end{bmatrix} = \begin{bmatrix} 0 & \frac{1}{2} & -\frac{1}{2} \\ \frac{1}{2} & -\frac{3}{4} & \frac{5}{4} \\ -\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \end{bmatrix} \begin{bmatrix} 6.9 \\ 7.2 \\ 3.0 \end{bmatrix} = \begin{bmatrix} 2.1 \\ 1.8 \\ -0.9 \end{bmatrix}$$

5) $\rho_1 = 0.9$, $\rho_2 = -0.9$, $\rho_3 = -0.9$,

$\rho_4 = 0.7$, $\rho_5 = -0.3$, $\rho_6 = 1.3$.

6) $|\rho_6| = \max_{1 \le i \le 6} |\rho_i|$; $\alpha = 6$.

7) $|\rho_6| > \rho$. $6 \notin I = [1,2,3]$.

8) $u = \text{sgn } \rho_6 = 1$.

9)

$$[\lambda_1, \lambda_2, \lambda_3] = [1,2,1] \begin{bmatrix} 0 & \frac{1}{2} & -\frac{1}{2} \\ \frac{1}{2} & -\frac{3}{4} & \frac{5}{4} \\ -\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \end{bmatrix}$$

$$= [\frac{1}{2}, -\frac{3}{4}, \frac{9}{4}] .$$

10) $\dfrac{u\ \lambda_3}{d_{33}} = 9 = \max_{1\le i \le 3} u\ \dfrac{\lambda_i}{d_{3i}}$; $k = 3$.

11)

$$D = \begin{bmatrix} \dfrac{1}{9} & \dfrac{1}{3} & -\dfrac{2}{9} \\[2mm] \dfrac{2}{9} & -\dfrac{1}{3} & \dfrac{5}{9} \\[2mm] -\dfrac{5}{9} & \dfrac{1}{3} & \dfrac{1}{9} \end{bmatrix}$$

12) $I = [1,2,6]$.

Step 2

4)

$$\begin{bmatrix} x_1 \\[2mm] x_2 \\[2mm] \rho \end{bmatrix} = \begin{bmatrix} \dfrac{1}{9} & \dfrac{1}{3} & -\dfrac{2}{9} \\[2mm] \dfrac{2}{9} & -\dfrac{1}{3} & \dfrac{5}{9} \\[2mm] -\dfrac{5}{9} & \dfrac{1}{3} & \dfrac{1}{9} \end{bmatrix} \begin{bmatrix} 6.9 \\[2mm] 7.2 \\[2mm] 7.0 \end{bmatrix} = \begin{bmatrix} \dfrac{29}{18} \\[2mm] \dfrac{136}{45} \\[2mm] -\dfrac{59}{90} \end{bmatrix}$$

5) $\rho_1 = \dfrac{59}{90}$, $\rho_2 = -\dfrac{59}{90}$, $\rho_3 = \dfrac{147}{90}$,

$\rho_4 = \dfrac{217}{90}$, $\rho_5 = -\dfrac{379}{90}$, $\rho_6 = -\dfrac{59}{90}$.

6) $|\rho_5| = \max_{1 \le i \le 6} |\rho_i|$. $\alpha = 5$.

7) $|\rho_5| > \rho$. $\quad 5 \notin I$.

8) $u = \text{sgn } \rho_5 = -1$.

9)
$$[\lambda_1, \lambda_2, \lambda_3] = [2, 4, -1] \begin{bmatrix} \frac{1}{9} & \frac{1}{3} & -\frac{2}{9} \\ \frac{2}{9} & -\frac{1}{3} & \frac{5}{9} \\ -\frac{5}{9} & \frac{1}{3} & \frac{1}{9} \end{bmatrix}$$

$$= [\frac{5}{3}, -1, \frac{5}{3}] .$$

10) $\dfrac{u \, \lambda_1}{d_{31}} = \dfrac{u \, \lambda_2}{d_{32}} = 3 = \max\limits_{1 \leq i \leq 3} u \, \dfrac{\lambda_i}{d_{3i}}$; $\quad k = 1 \quad$ or $\quad k = 2$.

We have shown in Chapter III the maximum of the above ratios
would be unique if Haar's condition were satisfied. We are
presented with two choices above, and we will consider
both of them.

Case 1 (k = 2)

11)
$$
D = \begin{bmatrix} \frac{2}{3} & -\frac{1}{3} & \frac{1}{3} \\ -\frac{1}{3} & \frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & \frac{2}{3} \end{bmatrix}
$$

12) I = [1,5,6] .

Step 3

4)
$$
\begin{bmatrix} x_1 \\ x_2 \\ \rho \end{bmatrix} = \begin{bmatrix} \frac{2}{3} & -\frac{1}{3} & \frac{1}{3} \\ -\frac{1}{3} & \frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & \frac{2}{3} \end{bmatrix} \begin{bmatrix} 6.9 \\ 11.1 \\ 7.0 \end{bmatrix} = \begin{bmatrix} \frac{97}{30} \\ \frac{14}{10} \\ \frac{29}{30} \end{bmatrix}
$$

5) $\rho_1 = -\frac{29}{30}$, $\rho_2 = -\frac{39}{10}$, $\rho_3 = -\frac{49}{30}$,

$\rho_4 = -\frac{25}{30}$, $\rho_5 = -\frac{29}{30}$, $\rho_6 = \frac{29}{30}$.

6) $|\rho_2| = \max_{1 \le i \le 6} |\rho_i|$; $\alpha = 2$.

7) $|\rho_2| > \rho$. $2 \notin I$.

8) $u = \text{sgn } \rho_2 = -1$.

9)
$$[\lambda_1, \lambda_2, \lambda_3] = [3,1,-1] \begin{bmatrix} \frac{2}{3} & -\frac{1}{3} & \frac{1}{3} \\ -\frac{1}{3} & \frac{1}{3} & 0 \\ 0 & -\frac{1}{3} & \frac{2}{3} \end{bmatrix}$$

$$= [\frac{5}{3}, -\frac{1}{3}, \frac{1}{3}] .$$

10) $\dfrac{u \; \lambda_1}{d_{31}} = \dfrac{5/3}{0}$, $\dfrac{u \; \lambda_2}{d_{32}} = -3$, $\dfrac{u \; \lambda_3}{d_{33}} = -1/2$.

Our computations break down at this point because of the division by zero. Let us therefore consider the second alternative at stage 10 in step 2.

Case 2   (k = 1)

11)

$$D = \begin{bmatrix} \frac{1}{15} & \frac{2}{5} & -\frac{1}{3} \\ \frac{2}{15} & -\frac{1}{5} & \frac{1}{3} \\ -\frac{1}{3} & 0 & \frac{2}{3} \end{bmatrix}$$

12)   I = [5,2,6] .

We now carry out the computations in Step 3. However, the computations will again break down at stage 10, for in calculating $u\,\lambda_2/d_{32}$, we will once more encounter a zero divisor. Thus the exchange algorithm has failed in this example with initial choice of I given by I = [1,2,3] .

There are two possible ways of working out of this dilemma. First, we can try to show in the context of the Dirichlet problem that Haar's condition is automatically satisfied. If this approach fails, we can try to modify the exchange algorithm so that it will work even though Haar's condition is violated. Let us examine the first of these alternatives.

4.2 Unisolvent Functions

Consider a set of functions $g_1,\ldots,g_n$ defined on a point set X. We then have

Definition 4.1. *The set of functions* $g_1,\ldots,g_n$ *is said to be unisolvent if the determinant*

$$(4.2) \qquad G = \begin{vmatrix} g_1(x_1) & \cdots & g_n(x_1) \\ \vdots & & \vdots \\ g_1(x_n) & \cdots & g_n(x_n) \end{vmatrix}$$

is non-zero for any n *distinct points* $x_1,\ldots,x_n$ *in* X.

Suppose now that we have a set of unisolvent functions $g_1,\ldots,g_n$ on a point set X. Then the coefficient matrix of the linear system

$$(4.3) \qquad \sum_{j=1}^{n} c_j g_j(x_i) = f(x_i) , \quad i=1,\ldots,m\geq n$$

obviously satisfies Haar's condition.

In the context of Dirichlet's problem, the linear system is

$$(4.4) \qquad \sum_{j=1}^{n} c_j g_j(x_i, y_i) = f(x_i, y_i) \ , \qquad i=1, \ldots, m > n \ ,$$

where the m points are selected from the contour C.
As a consequence of the following assertion, we should
always approximate f by at least three functions.

Theorem 4.1. *No set of two functions* g_1, g_2
are unisolvent on a simple closed curve C.

Proof: Suppose the determinant

$$(4.5) \qquad G = \begin{vmatrix} g_1(x_1, y_1) & g_2(x_1, y_1) \\ \\ g_1(x_2, y_2) & g_2(x_2, y_2) \end{vmatrix}$$

is non-zero for a certain pair of points (x_1, y_1) and
(x_2, y_2). By a continuous motion around the curve C,
these two points can be interchanged without becoming
coincident. In effect, we have then interchanged two
rows in the determinant in (4.5). Hence G has changed
sign. Therefore, at some intermediate position the
determinant must vanish, which was to be shown.

Let us now turn to the case where $g_1, \ldots, g_n$ are the
harmonic polynomials. We must then examine the behavior

of the following determinant:

(4.6)
$$G_n = \begin{vmatrix} 1 & \text{Re } z_1 & \text{Im } z_1 & \cdots & \text{Re } z_1^n & \text{Im } z_1^n \\ \cdot & \cdot & \cdot & \cdots & \cdot & \cdot \\ 1 & \text{Re } z_{2n+1} & \text{Im } z_{2n+1} & \cdots & \text{Re } z_{2n+1}^n & \text{Im } z_{2n+1}^n \end{vmatrix}$$

where $z_1, \ldots, z_{2n+1}$ are distinct points from C. To show a simple example of the vanishing of the above determinant, suppose $n = 1$. Then (4.6) reduces to

(4.7)
$$G_1 = \begin{vmatrix} 1 & x_1 & y_1 \\ 1 & x_2 & y_2 \\ 1 & x_3 & y_3 \end{vmatrix}$$

where $z_k = x_k + iy_k$, $k=1,2,3$. $G_1 = 0$ if the points (x_1,y_1), (x_2,y_2), (x_3,y_3) lie on any straight line. Therefore, if C contains a straight line segment, the functions 1, $\text{Re } z$, $\text{Im } z$ are not unisolvent on C.

The condition for the vanishing of G_n is that there exist some linear combination of the functions 1, $\text{Re } z$, $\text{Im } z$, $\ldots$, $\text{Re } z^n$, $\text{Im } z^n$, with coefficients not all zero, which vanishes at the points $z_1, \ldots, z_{2n+1}$. Therefore, the condition for the nonvanishing of G_n is

that the points $z_1, \ldots, z_{2n+1}$ do not lie on any curve of
the form $H_n(z) = 0$, where $H_n(z)$ is a linear combination
of the first $2n+1$ harmonic polynomials. If we wish this
to be true for any $2n+1$ points on C, then we must have
that C does not intersect any curve of the form $H_n(z) = 0$
at more than $2n$ points.

For an arbitrary simple closed curve C, this is a
difficult condition to check. Now $H_n(z) = 0$ is an
algebraic curve of degree n. An algebraic curve of degree
n can intersect an algebraic curve of degree 2 in at
most $2n$ distinct points. Hence, if C is an ellipse or
a circle, the harmonic polynomials are unisolvent on C.

A more extensive discussion of these matters can be
found in Curtiss [9,10].

4.3 Cheney's Perturbation Method

The results of the previous section are rather dis-
couraging in that the only interesting boundaries on which
we can be sure that the harmonic polynomials are unisolvent
are circles and ellipses. However, an alternative way of
handling the problem is a perturbation technique suggested
by Cheney ([13], p. 51). The method is suggested by the
following assertion.

Theorem 4.2. *If for a square matrix $A = (a_{ij})$,
$|A| = 0$, then for sufficiently small $\varepsilon \neq 0$, $|A - \varepsilon I| \neq 0$.*

Proof: Suppose that $|A - \lambda I| = 0$. Then λ is an eigen-value of the matrix (a_{ij}), $i,j = 1,\ldots,n$. There are at most n such eigenvalues, say $\lambda_1,\ldots,\lambda_n$. If ε is chosen such that $0 < |\varepsilon| < \max\limits_{1 \le j \le n} |\lambda_j|$, $\lambda_j \ne 0$, then the theorem is established.

To show how this method works, we again consider the problem of finding a minimax solution of (4.1). Since the vectors comprising the fifth and sixth equations are dependent, Theorem 4.1 suggests that we first determine the minimax solution to

$$2x_1 + x_2 = 6.9$$

$$3x_1 + x_2 = 7.2$$

$$x_1 + x_2 = 3.0$$

(4.9)

$$x_1 - x_2 = 1.0$$

$$(2+\varepsilon)x_1 + 4x_2 = 11.1$$

$$x_1 + (2+\varepsilon)x_2 = 7.0 \, ,$$

where $0 < |\varepsilon| < 4$. We will follow the same format as in Section 4.1. Initially, we again choose $I = [1,2,3]$.

Step 1

1) $T_1 = -2$, $T_2 = 1$, $T_3 = 1$.

2) $s_1 = -1$, $s_2 = 1$, $s_3 = 1$.

3)
$$D = \begin{bmatrix} 2 & 1 & -1 \\ 3 & 1 & 1 \\ 1 & 1 & 1 \end{bmatrix}^{-1} = \begin{bmatrix} 0 & \frac{1}{2} & -\frac{1}{2} \\ \frac{1}{2} & -\frac{3}{4} & \frac{5}{4} \\ -\frac{1}{2} & \frac{1}{4} & \frac{1}{4} \end{bmatrix} .$$

4)
$$\begin{bmatrix} x_1 \\ x_2 \\ \rho \end{bmatrix} = D \begin{bmatrix} 6.9 \\ 7.2 \\ 3.0 \end{bmatrix} = \begin{bmatrix} 2.1 \\ 1.8 \\ -0.9 \end{bmatrix} .$$

5) $\rho_1 = 0.9$, $\rho_2 = -0.9$, $\rho_3 = -0.9$,

   $\rho_4 = 0.7$, $\rho_5 \sim -0.3$, $\rho_6 \sim 1.3$.

In determining ρ_5 and ρ_6 , we have assumed that ε is

insignificantly small compared to 0.3. We have replaced
the equality sign by the symbol "~" whenever a similar
process has been carried out.

6) $\quad |\rho_6| = \max_{1 \le i \le 6} |\rho_i| \; ; \quad \alpha = 6$.

7) $\quad |\rho_6| > \rho$. $6 \notin I$.

8) $\quad u = 1$.

9) $\quad [\lambda_1, \lambda_2, \lambda_3] = [1, 2+\varepsilon, 1] \; D = [\frac{1}{2}(1+\varepsilon), \; -\frac{3}{4}(1+\varepsilon), \; \frac{1}{4}(9+5\varepsilon)]$.

10) $\quad \dfrac{u \, \lambda_3}{d_{33}} \sim 9 = \max_{1 \le i \le 3} u \, \dfrac{\lambda_i}{d_{3i}} \; ; \quad k = 3$.

11) $\quad D = \dfrac{1}{9+5\varepsilon} \begin{bmatrix} 1+\varepsilon & 3+\varepsilon & -2 \\ 3 & -3 & 5 \\ -5-3\varepsilon & 3+2\varepsilon & 1 \end{bmatrix}$.

12) $I = [1,2,6]$.

Step 2

4)
$$\begin{bmatrix} x_1 \\ \\ x_2 \\ \\ \rho \end{bmatrix} = \frac{1}{10(9+5\epsilon)} \begin{bmatrix} 145 + 141\epsilon \\ \\ 272 \\ \\ -59 - 63\epsilon \end{bmatrix}$$

5) $\rho_1 \sim \frac{59}{90}$, $\rho_2 \sim -\frac{59}{90}$, $\rho_3 \sim -\frac{147}{90}$,

$\rho_4 \sim \frac{217}{90}$, $\rho_5 \sim -\frac{379}{90}$, $\rho_6 \sim -\frac{59}{90}$.

6) $|\rho_5| = \max_{1 \leq i \leq 6} |\rho_i|$; $\alpha = 5$.

7) $|\rho_5| > \rho$. $5 \notin I$.

8) $u = -1$.

9) $[\lambda_1, \lambda_2, \lambda_3] = [2+\varepsilon, 4, -1]$ D

$$= \frac{1}{9+5\varepsilon} [15+6\varepsilon+\varepsilon^2, \ -9+3\varepsilon+\varepsilon^2, \ 15-2\varepsilon] \ .$$

10) $\dfrac{u \ \lambda_1}{d_{31}} = \dfrac{15+6\varepsilon+\varepsilon^2}{5+3\varepsilon}$; $\dfrac{u \ \lambda_2}{d_{32}} = \dfrac{9-3\varepsilon-\varepsilon^2}{3+2\varepsilon}$; $\dfrac{u \ \lambda_3}{d_{33}} = -15+2\varepsilon$.

For $\varepsilon > 0$ and sufficiently small, the maximum of the above

ratios is $\dfrac{u \ \lambda_1}{d_{31}}$. Therefore, $k = 1$.

11) $D = \dfrac{1}{(9+5\varepsilon)(15+6\varepsilon+\varepsilon^2)} \times$

$$\begin{bmatrix} (1+\varepsilon)(9+5\varepsilon) & 54+39\varepsilon+5\varepsilon^2 & -45-25\varepsilon \\ 2(9+5\varepsilon) & -27-24\varepsilon-5\varepsilon^2 & 45+34\varepsilon+5\varepsilon^2 \\ (-5-3\varepsilon)(9+5\varepsilon) & 36\varepsilon+29\varepsilon^2+5\varepsilon^2 & 90+41\varepsilon-5\varepsilon^2 \end{bmatrix}$$

12) $I = [5, 2, 6]$.

<u>Step 3</u>

4)
$$
\begin{bmatrix} x_1 \\ x_2 \\ \rho \end{bmatrix} = \frac{1}{10(9+5\varepsilon)(15+6\varepsilon+\varepsilon^2)} \begin{bmatrix} 1737 + 2612\varepsilon + 915\varepsilon^2 \\ 3204 + 1762\varepsilon - 10\varepsilon^2 \\ 1305 - 310\varepsilon - 7\varepsilon^2 + 360\varepsilon^3 \end{bmatrix}
$$

5) $\rho_1 \sim \frac{2637}{1350}$, $\rho_2 \sim \frac{1305}{1350}$, $\rho_3 \sim -\frac{891}{1350}$,

$\rho_4 \sim \frac{2817}{1350}$, $\rho_5 \sim -\frac{1305}{1350}$, $\rho_6 \sim \frac{1305}{1350}$.

6) $|\rho_4| = \max_{1\leq i \leq 6} |\rho_i|$; $\alpha = 4$.

7) $|\rho_4| > \rho$. $4 \notin I$.

8) $u = 1$.

9) $[\lambda_1,\lambda_2,\lambda_3] = [1,-1,1]$ D

$$= \frac{1}{(9+5\varepsilon)(15+6\varepsilon+\varepsilon^2)} \quad \times$$

$$[(-6-2\varepsilon)(9+5\varepsilon), \ 81+99\varepsilon+39\varepsilon^2+5\varepsilon^3, -18\varepsilon-10\varepsilon^2].$$

10) $\displaystyle \frac{u\,\lambda_2}{d_{32}} = \max_{1\le i\le 3} u\,\frac{\lambda_i}{d_{3i}} \ ; \quad k = 2$.

11) $D = \begin{bmatrix} \dfrac{1}{3+\varepsilon} & \dfrac{6+\varepsilon}{(3+\varepsilon)^2} & -\dfrac{3}{(3+\varepsilon)^2} \\[4mm] 0 & -\dfrac{1}{3+\varepsilon} & \dfrac{1}{3+\varepsilon} \\[4mm] -\dfrac{1}{3+\varepsilon} & \dfrac{4\varepsilon+\varepsilon^2}{(3+\varepsilon)^2} & \dfrac{6+\varepsilon}{(3+\varepsilon)^2} \end{bmatrix}$

12) $I = [5,4,6]$.

Step 4

4) $$\begin{bmatrix} x_1 \\ \\ x_2 \\ \\ \rho \end{bmatrix} = \begin{bmatrix} \dfrac{183+121\varepsilon}{10(3+\varepsilon)^2} \\ \\ \dfrac{6}{3+\varepsilon} \\ \\ \dfrac{87-\varepsilon+10\varepsilon^2}{10(3+\varepsilon)^2} \end{bmatrix}$$

5) $\rho_1 \sim \dfrac{75}{90}$, $\rho_2 \sim -\dfrac{81}{90}$, $\rho_3 \sim -\dfrac{93}{90}$,

$\rho_4 \sim \dfrac{87}{90}$, $\rho_5 \sim -\dfrac{87}{90}$, $\rho_6 \sim \dfrac{87}{90}$.

6) $|\rho_3| = \max_{1 \le i \le 6} |\rho_i|$; $\alpha = 3$.

7) $|\rho_3| > \rho$. $3 \notin I$.

8) $u = -1$.

9) $[\lambda_1, \lambda_2, \lambda_3] = [1,1,-1]\, D$

$$= [\ \frac{2}{3+\epsilon}\ ,\ \frac{3-4\epsilon-\epsilon^2}{(3+\epsilon)^2}\ ,\ -\frac{6}{(3+\epsilon)^2}\ ]\ .$$

10) $2 \sim \dfrac{u\,\lambda_1}{d_{31}} = \max_{1\le i\le 3} u\,\dfrac{\lambda_i}{d_{3i}}\ ;\quad k = 1\ .$

11) $D = \begin{bmatrix} \dfrac{1}{2} & \dfrac{1}{2} & 0 \\[2ex] 0 & -\dfrac{1}{3+\epsilon} & \dfrac{1}{3+\epsilon} \\[2ex] -\dfrac{1}{2} & \dfrac{1+\epsilon}{6+2\epsilon} & \dfrac{1}{3+\epsilon} \end{bmatrix}$

12) $I = [3,4,6]\ .$

Step 5

4) $\begin{bmatrix} x_1 \\[1ex] x_2 \\[1ex] \rho \end{bmatrix} = \begin{bmatrix} 2 \\[1ex] \dfrac{6}{3+\epsilon} \\[1ex] \dfrac{6-2\epsilon}{6+2\epsilon} \end{bmatrix}$

5) $\rho_1 \sim 0.9$, $\rho_2 \sim -0.8$, $\rho_3 \sim -1.0$,

   $\rho_4 \sim 1.0$, $\rho_5 \sim -0.9$, $\rho_6 \sim 1.0$.

6) $|\rho_3| = |\rho_4| = |\rho_6| = \max_{1 \le i \le 6} |\rho_i|$; $\alpha = 3$ or 4 or 6 .

7) $|\rho_3| = |\rho_4| = |\rho_6| > \rho$; $\alpha \in I$.

The last line shows that the critical subsystem has been reached. To find the minimax solution of (4.1) we let $\epsilon \to 0$ to obtain $x_1 = x_2 = 2$ and $\rho = 1$.

The technique shown above is not convenient for use on a computer since ϵ has been left as a variable in the computations. The same method employed on a computer would require the use of symbol manipulation, which is not readily available. The only recourse would seem to be to introduce a sequence of ϵ's say $\epsilon_1 > \epsilon_2 > \ldots \epsilon_n \to 0$ and examine what happens. If no significant changes are observed after a certain ϵ_N , then one could accept the minimax solution to the problem with ϵ_N appearing as a good approximation to the minimax solution of the original problem.

The problem of finding the minimax solution of (4.1) was
attacked using the computer program MINMAX given in the
Appendix. This program assumes that Haar's condition is
satisfied. The correct solution was obtained using both
single and double precision arithmetic, using as initial
guess I = [1,2,3] . The computations did not break down
(as in the exact calculation shown in Section 4.1) because
roundoff errors introduced during computations yielded very
small numbers rather than zeros at the critical stages.
In effect, the roundoff errors perturbed the original
system, which did not satisfy Haar's condition, into an
adjacent system that did. We present some of the results
of the single precision computations to 3 significant
digits. (The initial subsystem is given by I = [1,2,3]).

Step 2

1) $T_1 = -2$, $T_2 = 1$, $T_3 = 1$.

2) $s_1 = -1$, $s_2 = 1$, $s_3 = 1$.

3)
$$D = \begin{bmatrix} 0.0 & 0.5 & -0.5 \\ 0.5 & -0.75 & 1.25 \\ -0.5 & 0.25 & 0.25 \end{bmatrix}$$

4)
$$\begin{bmatrix} x_1 \\ x_2 \\ \rho \end{bmatrix} = \begin{bmatrix} 2.10 \\ 1.80 \\ -0.90 \end{bmatrix}$$

5) $\rho_1 = 0.9$, $\rho_2 = -0.9$, $\rho_3 = -0.9$,

$\rho_4 = 0.7$, $\rho_5 = -0.3$, $\rho_6 = 1.30$.

6) $|\rho_6| = \max_{1 \leq i \leq 6} |\rho_i|$; $\alpha = 6$.

7) $|\rho_6| > \rho$. $6 \notin I$.

8) $u = 1$.

9) $[\lambda_1,\lambda_2,\lambda_3] = [0.5,-0.75,2.25]$.

10) $9.0 = \dfrac{u \; \lambda_3}{d_{33}} = \max_{1\le i \le 3} \; u \; \dfrac{\lambda_i}{d_{3i}} \; ; \quad k = 3$.

11)
$$D = \begin{bmatrix} 0.111 & 0.333 & -0.222 \\ 0.222 & -0.333 & 0.556 \\ -0.556 & 0.333 & 0.111 \end{bmatrix}$$

12) $I = [1,2,6]$.

Step 2

4)
$$\begin{bmatrix} x_1 \\ x_2 \\ \rho \end{bmatrix} = \begin{bmatrix} 1.61 \\ 3.02 \\ -0.656 \end{bmatrix}$$

5) $\rho_1 = 0.656$, $\rho_2 = -0.656$, $\rho_3 = -1.63$,

$\rho_4 = 2.41$, $\rho_5 = -4.21$, $\rho_6 = -0.656$.

6) $|\rho_5| = \max\limits_{1 \le i \le 6} |\rho_i|$; $\alpha = 5$.

7) $|\rho_5| > \rho$. $5 \notin I$.

8) $u = -1$.

9) $[\lambda_1, \lambda_2, \lambda_3] = [1.67, -1.0, 1.67]$.

10) $3.0 = \dfrac{u \lambda_2}{d_{32}} = \max\limits_{1 \le i \le 3} u \dfrac{\lambda_i}{d_{3i}}$; $k = 2$.

11)
$$D = \begin{bmatrix} 0.667 & -0.333 & 0.333 \\ -0.333 & 0.333 & 0.212 \times 10^{-6} \\ -0.212 \times 10^{-6} & -0.333 & 0.667 \end{bmatrix}$$

12) $I = [1, 5, 6]$.

Note in stage (11) that $d_{31} \neq 0$. In the exact computation, $d_{31} = 0$. The computations therefore do not break down. The algorithm then considers successively the subsystems [2,5,6], [3,5,6], and [3,4,6]. It reports that the last subsystem is the critical one and gives the minimax solution as $x_1 = x_2 = 2.0$ and the associated minimax error as $\rho = 1.0$.

This result is somewhat encouraging. It should be noted, however, that the reverse process is also possible: that is, it is possible for roundoff errors to perturb a system satisfying Haar's condition into an adjacent one that violates Haar's condition. Individual cases will probably require individual treatment.

CHAPTER V

MINIMAX APPROXIMATIONS ON A CONTINUUM

5.1 Introduction

In the previous two chapters, we solved the problem of determining a vector $c = (c_1, \ldots, c_n)$ to minimize

$$(5.1) \qquad M_1 = \max_{1 \leq i \leq m} \left| f(x_i, y_i) - \sum_{j=1}^{n} c_j g_j(x_i, y_i) \right| ,$$

where (x_i, y_i), $i = 1, \ldots, m$, were distinct points on the contour C. Intuitively, we feel that if the number of points becomes large, leaving no wide gaps along C, the minimax error associated with minimizing (5.1) should be a good approximation to the minimax error associated with minimizing

$$(5.2) \qquad M_2 = \max_{(x,y) \varepsilon C} \left| f(x,y) - \sum_{j=1}^{n} c_j g_j(x,y) \right| .$$

Thus, one approach to finding the minimax error on the curve C would be as follows: Compute the minimax error for m points. Double the number of points and again

calculate the minimax error. If the minimax error does
not change significantly, then stop. Otherwise, once more
double the number of points and continue as before.

A more elegant approach was first conceived by Remes
in 1934. The statement of the algorithm and the proof
of its convergence is from Cheney ([3], pp. 95-96).

5.2 Remes Algorithm

This algorithm only requires that the functions
$f, g_1, \ldots, g_n$ be continuous on the curve C . We wish to
determine a vector $c = (c_1, \ldots, c_n)$ for which the deviation

$$(5.3) \qquad M(c) = \max_{s \varepsilon C} \left| f(s) - \sum_{j=1}^{n} c_j g_j(s) \right|$$

is a minimum. We may assume without loss of generality
that $g_1, \ldots, g_n$ are linearly independent, for if they are
not, we can replace them by a smaller number of functions
that are independent without raising the minimum deviation

$$(5.4) \qquad \rho = \min_{c} M(c) .$$

We define a residual function $r(c,s)$ by

$$(5.5) \qquad r(c,s) = f(s) - \sum_{j=1}^{n} c_j g_j(s) \; .$$

Then Remes algorithm is as follows: At the kth step we are given a finite subset S^k of C . Select a coefficient vector c^k to minimize the function

$$(5.6) \qquad M^k(c) = \max_{s \varepsilon S^k} |r(c,s)| \; .$$

This can be accomplished by the procedure given in Chapter III. Choose $s^k \varepsilon C$ to maximize $|r(c^k,s)|$. Thus,

$$(5.7) \qquad |r(c^k,s^k)| = M(c^k) \; .$$

Now start again with the finite set $S^{k+1} = S^k \cup \{s^k\}$. At the beginning, S^1 may be arbitrary except that the set of n-tuples $\hat{s} = [g_1(s),\ldots,g_n(s)]$ corresponding to $s \varepsilon S^1$ should be of rank n . The algorithm stops when $|r(c^k,s^k)|$ is insignificantly different from $\min_{c} \max_{s \varepsilon S^k}$ $|r(c,s)|$. A theorem associated with Remes algorithm is

Theorem 5.1. $\lim_{k \to \infty} M^k(c^k) = \rho$. *The sequence* c^k *is bounded and its cluster points minimize* M .

Proof: We define

(5.8)
$$\|c\|_1 = \sum_{j=1}^{n} |c_j| .$$

It can be shown that this function defines a norm in the linear space R_n . By our assumption on S^1 , it follows that

(5.9)
$$\phi = \min_{\|c\|_1 = 1} \max_{s \varepsilon S^1} | \sum_{j=1}^{n} c_j g_j(s) | > 0 .$$

Consequently,

$$M^1(c) = \max_{s \varepsilon S^1} |f(s) - \sum_{j=1}^{n} c_j g_j(s)|$$

$$= \| f - \sum_{j=1}^{n} c_j g_j \|$$

(5.10)

$$\geq \| \sum_{j=1}^{n} c_j g_j \| - \| f \|$$

$$= \max_{s \varepsilon S^1} | \sum_{j=1}^{n} c_j g_j(s) | - \max_{s \varepsilon S^1} |f(s)| .$$

Now let

(5.11) $$\|c\|_1 = K .$$

Then

(5.12) $$\max_{s \varepsilon S^1} | \sum_{j=1}^{n} c_j g_j(s) | = K \max_{s \varepsilon S^1} | \sum_{j=1}^{n} c_j/K \, g_j(s) | .$$

Therefore, we may rewrite (5.10) as

(5.13) $$M^1(c) \geq \|c\|_1 \max_{s \varepsilon S^1} | \sum_{j=1}^{n} c_j g_j(s) | - \max_{s \varepsilon S^1} |f(s)| ,$$

and hence by (5.9),

(5.14) $$M^1(c) \geq \|c\|_1 \phi - \max_{s \varepsilon S^1} |f(s)| .$$

Now suppose that $\|c\|_1 > 2 \max_{s \varepsilon S^1} |f(s)|/\phi$. Then

$$M^k(c) \geq M^1(c)$$

(5.15)
$$> \max_{s \in S^1} |f(s)|$$

$$\geq M^k(\Theta) \, ,$$

where Θ is the zero vector. Therefore, c does not enter the competition to minimize any of the functions M^k . Thus the sequence c^k generated by the algorithm is bounded. Now we have

(5.16)
$$S^k \subset S^{k+1} \subset C \, ,$$

and hence

(5.17)
$$M^k(c) \leq M^{k+1}(c) \leq M(c)$$

for any vector c . Therefore,

(5.18)
$$M^k(c^k) \leq M^{k+1}(c^{k+1}) \leq \rho \, .$$

Thus for some $\epsilon \geq 0$, $M^k(c^k) \rightarrow \rho - \epsilon$. We must show that $\epsilon = 0$. Since

$$(5.19) \qquad |r(b,s) - r(c,s)| = |\sum_{j=1}^{n} (c_j - b_j)\, g_j(s)|$$

$$\leq T \|b-c\|_1 \, ,$$

where

$$(5.20) \qquad T = \max_{1 \leq j \leq n} \max_{s \in C} |g_j(s)| \, ,$$

it follows that

$$(5.21) \qquad |r(b,s)| \leq |r(c,s)| + T \|b-c\|_1$$

and furthermore

$$(5.22) \qquad M(b) \leq M(c) + T \|b-c\|_1 \, .$$

Suppose now that $\epsilon > 0$. Let b denote any cluster point

of the sequence c^k . For any $\delta > 0$, we may find an index k such that

$$(5.23) \qquad \|b-c^k\|_1 < \delta$$

and an index $i > k$ such that

$$(5.24) \qquad \|b-c^i\|_1 < \delta \ .$$

Then

$$(5.25) \qquad \|c^i-c^k\|_1 \leq 2\delta$$

and

$$\rho \leq M(b) \leq M(c^k) + T\delta$$

$$= |r(c^k,s^k)| + T\delta$$

$$(5.26) \qquad \leq |r(c^i,s^k)| + 3T\delta$$

$$\leq M^i(c^i) + 3T\delta$$

$$\leq \rho - \epsilon + 3T\delta \ .$$

Picking δ small enough so that $3T\delta < \varepsilon$ yields a

contradiction. Therefore, $\varepsilon = 0$ and (5.26) can be

rewritten

(5.27) $\rho \leq M(b) \leq \rho + 3T\delta$.

Letting $\delta \to 0$ proves the assertion.

 Remes algorithm requires that we locate the extremum

of the residual function $r(c,s)$. This can be a difficult

problem even in well-behaved cases. There would be no

problem in constructing an example to defeat any particular

method.

 We have chosen the obvious and simple method of

calculating the residual on a discrete set of points chosen

in advance. Originally, this discrete set was chosen as

a set of equally spaced points. Examination of the residual

function, however, has shown that it has approximately the

form of a Chebyshev polynomial ([13], p. 60), with the

extrema crowded towards the ends of the range under

consideration. Therefore, the discrete point set was chosen

to be the zeros of a Chebyshev polynomial. As a rule of

thumb, if we are approximating f by n functions

$g_1, \ldots, g_n$, then the discrete point set was chosen to be

the zeros of the (20n)th Chebyshev polynomial, although

no significant change in the minimax error was evident if

we took the discrete point set to be the zeros of the
Chebyshev polynomial of (10n)th degree or higher.

CHAPTER VI

NUMERICAL RESULTS

The present experiments were performed using the IBM System 360 FORTRAN IV source language and the IBM 360/67 computer system at the University of Alberta, Edmonton.

6.1 Test Problems and Results

For clarity, we shall elaborate the meaning of several terms which occur in the tabulation of the data.

Boundary C: The boundary of the region of space over which the solution is sought.

Boundary values: A continuous function f defined on the boundary. If there are more than one, then the functions will be labelled (a), (b), ...

Approximating functions: Harmonic functions selected from a set of harmonic functions which is complete over the space of functions which are continuous on the boundary. Usually, these functions will be selected from the set of harmonic polynomials (see Section 2.3).

Boundary-value problem: This will be either the Dirichlet problem or the mixed problem for Laplace's equation (see [16], p. 169, for the definition of the mixed problem).

Minimax error, ρ_n : The computed minimax error when n approximating functions are employed.

Minimax solution: This will be written in the form $u \sim c_1 g_1 + \ldots + c_n g_n$, where u is the required solution of the Dirichlet problem, $g_1, \ldots, g_n$ are the approximating functions, and $c_1, \ldots, c_n$ are the computed coefficients that minimize the maximum error on the boundary.

Least-Squares solution: This will be of the same form as the minimax solution except that the coefficients minimize the sum of the squares of the errors at discrete points along the boundary.

Error at isolated interior point(s): The absolute value of the difference between the theoretical value and the computed value at the given point(s). This can be given only if the true solution is known.

6.1.1 Problem 1 For complete details see Davis [12].

Boundary: A bean-shaped region was obtained from a free-hand drawing on coordinate paper (see Figure 6.1). The boundary is "defined" by means of 43 points (see Table 6.1). These points are not distributed equally on the boundary; more points appear where the curvature is greater.

Figure 6.1 Bean-shaped region of Problem 1

Table 6.1 Points defining the bean-shaped region

Pt. no.	Abscissa	Ordinate
1	.000	.110
2	-.050	.108
3	-.100	.115
4	-.160	.150
5	-.220	.205
6	-.320	.300
7	-.400	.358
8	-.500	.420
9	-.550	.436
10	-.600	.430
11	-.644	.400
12	-.660	.350
13	-.655	.300
14	-.635	.200
15	-.595	.100
16	-.552	.000
17	-.500	-.105
18	-.440	-.200
19	-.400	-.250
20	-.350	-.300
21	-.300	-.344
22	-.204	-.400
23	-.100	-.436
24	.000	-.448
25	.100	-.442

Table 6.1 (Continued)

Pt. no.	Abscissa	Ordinate
26	.230	-.400
27	.300	-.350
28	.353	-.300
29	.430	-.200
30	.477	-.100
31	.510	.000
32	.522	.100
33	.520	.160
34	.500	.240
35	.456	.300
36	.400	.330
37	.360	.337
38	.300	.320
39	.250	.290
40	.200	.245
41	.150	.200
42	.100	.160
43	.050	.128

Boundary values:

(a) $f = 0$, $x \leq 0$; $f = x$, $x > 0$.

(b) $f = 0$, $x \leq 0$; $f = x^2$, $x > 0$.

(c) $f = 0$, $x \leq 0$; $f = x^3$, $x > 0$.

(d) $f = 0$, $x \leq 0$; $f = x^4$, $x > 0$.

(e) $f = 0$, $x \leq 0$; $f = x^5$, $x > 0$.

(f) $f = x^2 + y^2$ (Torsion problem).

(g) $f = \exp(x^2 - xy + 2y^2)$.

(h) $f = \ln[(x + 1.5)^2 + y^2]$.

(i) $f = \ln[x + y + (y-1)^2]$.

(j) $f = e^x \cos y + \ln[x^2 + (y-1)^2]$.

Approximating functions: 1, $\mathrm{Re}\, z^n$, $\mathrm{Im}\, z^n$; $n=1,2,\ldots$.

Boundary value problem: The Dirichlet problem.

Minimax error: See Table 6.2.

<u>Minimax solution</u>: This is for function (j) for n = 5 .

$$u \sim 1.00244 + .997007 \text{ Re } z - 1.99034 \text{ Im } z$$

$$+ 1.47952 \text{ Re } z^2 - .0110482 \text{ Im } z^2$$

$$+ .187437 \text{ Re } z^3 + .628927 \text{ Im } z^3$$

$$- .345989 \text{ Re } z^4 + .0152750 \text{ Im } z^4$$

$$- .133709 \text{ Re } z^5 - .322899 \text{ Im } z^5$$

<u>Least squares solution</u>: The coefficients are quoted from
[12]. This is for the same problem as above.

$$u \sim 1.00173 + .997339 \text{ Re } z - 1.99119 \text{ Im } z$$

$$+ 1.48065 \text{ Re } z^2 - .00949996 \text{ Im } z^2$$

$$+ .188957 \text{ Re } z^3 + .623677 \text{ Im } z^3$$

$$- .355600 \text{ Re } z^4 + .024526 \text{ Im } z^4$$

$$- .11960 \text{ Re } z^5 - .28034 \text{ Im } z^5$$

For this approximation the maximum error on
the boundary is .0069.

Table 6.2 Results for bean-shaped region

n	Minimax error ρ_n for function
	(a)	(b)	(c)	(d)	(e)
1	.142	.0823	.0477	.0266	.0145
2	.0870	.0410	.0232	.0135	.00776
3	.0725	.0284	.0131	.00657	.00357
4	.0613	.0234	.00961	.00378	.00156
5	.0562	.0182	.00757	.00321	.00138
6	.0480	.0164	.00662	.00276	.00123
7	.0446	.0141	.00546	.00233	.00105
8	.0414	.0126	.00483	.00202	.000909
9	.0382	.0113	.00416	.00174	.000783

n	Minimax error ρ_n for function
	(f)	(g)	(h)	(i)	(j)
1	.190	.544	.613	.384	.277
2	.148	.338	.0322	.192	.0858
3	.0878	.184	.00930	.0940	.0258
4	.0707	.116	.00250	.0429	.0110
5	.0599	.0964	.000759	.0216	.00434
6	.0505	.0810	.000226	.0122	.00179
7	.0433	.0698	.0000707	.00879	.000738
8	.0371	.0607	.0000204	.00645	.000321
9	.0298	.0487	.000000545	.00491	.000134

<u>Error at isolated interior points</u>: The results (see Table

6.3) pertain to function (j) for both the

minimax and least-squares solution.

Table 6.3 Discrepancies at points interior

to bean along y = 0

x	minimax discrepancies	least-squares discrepancies
-.5	.0033	.0014
-.4	.0021	.0010
-.3	.0020	.0011
-.2	.0023	.0015
-.1	.0025	.0017
0	.0024	.0017
.1	.0020	.0013
.2	.0013	.0007
.3	.0006	.0001
.4	.0005	.0009

6.1.2 <u>Problem 2</u>

<u>Boundary</u>: The ellipse $x^2 + 4y^2 = 1$.

Boundary values:

    (a) $f = 0$, $x \leq 0$; $f = x$, $x > 0$.

    (b) $f = 0$, $x \leq 0$; $f = x^2$, $x > 0$.

    (c) $f = 0$, $x \leq 0$; $f = x^3$, $x > 0$.

    (d) $f = \ln[(x + 1.5)^2 + y^2]$.

Approximating functions: 1, Re z^n; $n = 1, 2, \ldots$. Only the real parts of z^n contribute since the boundary and boundary functions are symmetric with respect to the x axis. Thus only the top half of the boundary need be considered.

Boundary-value problem: The Dirichlet problem.

Minimax error: See Table 6.4.

Table 6.4   Results for the ellipse

n	Minimax error ρ_n for function			
	(a)	(b)	(c)	(d)
1	.250	.281	.317	.385
2	.0577	.0856	.134	.139
3	.0525	.0138	.0370	.0577
4	.0311	.0138	.00442	.0258
5	.0308	.0531	.00444	.0120
6	.0210	.0531	.00133	.00571
7	.0209	.00277	.00133	.00278
8	.0159	.00277	.000566	.00137

Exact value of solution for function (d)   at the point
(0,0) :   u(0,0) = ln 2.25 = 0.8109 .

Computed value of approximating solution at the point
(0,0)   for n approximating functions: See Table 6.5.

Table 6.5   Computed solution for function

(d)   at   (0,0)

n	computed value at (0,0)	discrepancy between computed and true solution
2	.6105	.2004
3	.8544	.0435
4	.8207	.0098
5	.8077	.0032
6	.8103	.0006
7	.8112	.0003
8	.8110	.0001
9	.8109	.0000

6.1.3   Problem 3

Boundary: Square given by the lines $x = 0$, $x = 1$, $y = 0$,

y = 1 .

Boundary values: (a)

$$
f = \begin{cases} 1, & y = 1, \ 0 < x < 1 . \\ \\ 0, & \text{otherwise} . \end{cases}
$$

(b)

$$
f = \begin{cases} x^2 - x, & y = 1, \ 0 < x < 1 . \\ \\ 0, & \text{otherwise} . \end{cases}
$$

Approximating functions: sinh πny sin πnx, n=1,2,... .

These functions are harmonic and satisfy the
homogeneous boundary conditions. Because of
the symmetry of the problem along the line
$x = \frac{1}{2}$, the approximating functions used are
sinh $\pi(2n-1)y$ sin $\pi(2n-1)x$, n=1,2,... .

Minimax error: See Table 6.6.

Table 6.6 Results for Square

No. of functions	Minimax error ρ_n for function	
	(a)	(b)
2	1.00	.0024
3	1.00	.0010
4	1.00	.00054
5	1.00	.00034
6	1.00	.00023
7	1.00	.00016
8	1.00	.00012
9	1.00	.00010

Exact solution for function (a) at (0.5,0.5): u(0.5,0.5) =

0.25 (see [22], p. 105).

Computed value of solution at (0.5,0.5) for n

approximating functions: See Table 6.7.

Table 6.7 Computed solution for function

(a) at (0.5,0.5)

n	value
2	.2474
3	.2519
4	.2505
5	.2504
6	.2504
7	.2504
8	.2455
9	.2501

6.1.4 __Problem 4__ For complete details, see Charmonman [2].

__Boundary__: Rectangle given by the lines x = 0, x = -2,

y = 0, y = -1 .

__Boundary values__:

$$f = \begin{cases} x, \ y = -1, \ -1.5 < x < 0 \ . \\ \\ 0, \ \text{otherwise} \ . \end{cases}$$

Boundary-value problem: The mixed problem. The normal

　　　　derivative of the solution is specified along

　　　　y = -1, -2 < x < -1.5 . Elsewhere, the value

　　　　of the solution is given.

Approximating functions:　i)　$\sinh \frac{\pi n y}{-2} \sin \frac{\pi n x}{-2}$, n=1,2,...

　　　　　　　　　　　　where the value of the

　　　　　　　　　　　　solution is given.

　　　　　　　　ii)　$\frac{\pi n}{-2} \cosh \frac{\pi n y}{-2} \sin \frac{\pi n x}{-2}$, n=1,2,...

　　　　　　　　　　　　where the normal derivative

　　　　　　　　　　　　of the solution is specified.

Minimax error: See Table 6.8.

Table 6.8　Results for Rectangle

No. functions	Minimax error ρ
2	.7461
3	.6317
4	.6140
5	.5875
6	.5264
7	.5204
8	.5052
9	.4612
10	.4598

6.2 Comments on Numerical Results

Problem 1: Davis [12] mentions this problem as
a difficult test case because of the nonconvexity of the
domain. Nevertheless, we were interested in testing our
procedure for a fairly intricate region.

Boundary functions (a) to (e) yield a family of
boundary values of increasing smoothness. The table of
minimax errors show a decrease in the maximum boundary
error as the boundary functions become smoother. In
addition, the rate of decrease of the minimax error as
we introduce more approximating functions is greater for
smoother boundary functions.

Function (j) is harmonic over the region over which
the solution is required and hence the boundary function
is the solution to the Dirichlet problem. This test case
was tried in [12] using a least-squares approach mentioned
in Chapter I. We note that the maximum error along the
boundary is smaller using our minimax procedure than the
maximum error given by Davis' least-squares technique.
However, the discrepancies in the interior seem to be
larger for the minimax technique. In both approaches, the
discrepancies in the interior are less than the computed
maximum error along the boundary.

Problem 2:   For test functions we selected four of the functions which were tried in Problem 1.   (Boundary functions (a), (b), (c), (d) for the ellipse correspond to boundary functions (a), (b), (c), (h) for the bean.) For the family of functions (a) to (c) , we see that the minimax errors for the ellipse are less than the corresponding minimax errors for the bean.   This is probably due to the nonconvexity of the bean-shaped region.   On the other hand, the minimax errors for function (d) are bigger for the ellipse.   This is a result of the singularity of the boundary function being closer to the ellipse than to the bean.   For function (d) , we have listed the differences, for increasing numbers of approximating functions, between the true solution and the computed approximate solution at an interior point.   The differences become smaller as the number of approximating functions increases.   We note also that the discrepancies are well within the computed maximum error along the boundary.

Problem 3:   The two boundary functions used show the difference in results between a case where the boundary function possesses a discontinuity on the boundary (function (a)) and a case where it does not (function (b)). The solution for the Dirichlet problem using function (a) can be written as an infinite series whose value is known at the interior point (0.5,0.5) .   We have shown how the

computed value approaches the exact value at this point as more approximating functions are introduced. The approach to the true solution is quite uniform except for 8 approximating functions where the sudden divergence is something of a mystery, although it in no way contradicts any theory that we have presented.

Problem 4: Although we have not dealt with the mixed problem in the text of this thesis, we were interested to see how our minimax procedure would fare. For approximating the normal derivative of the solution along the boundary we have used a standard technique in numerical analysis: If L designates a linear operator and u a function, then we approximate L(u) by L(approximation to u). It is interesting to note that the minimax error still decreases as the number of approximating functions increases. This leads us to believe that the technique may be used with success in other problem areas, such as the Neumann problem for Laplace's equation and problems involving multiply-connected domains. The large minimax errors are probably due to the singularity at the point $(-1.5, -1)$.

CHAPTER VII

CONCLUSIONS AND SUGGESTIONS

FOR FURTHER RESEARCH

7.1 Conclusions

The following conclusions have been drawn from the numerical computations:

1) As opposed to finite-difference techniques, the method presented is easy to apply to boundaries of irregular shape (for example, the bean-shaped region). Further-more, obtaining the approximate solution at any desired point(s) in the interior of the region is relatively simple.

2) The technique is certainly as good as other collocation techniques such as the point-matching method and the least-squares method. The maximum error is usually a strictly decreasing function of the number of approxi-mating harmonic functions. This is a virtue that the point-matching method does not always possess (see Section 1.3.1). The maximum error along the boundary is smaller for this method than for the least-squares method, and hence the bound of the error between the true solution and the computed solution is smaller. Furthermore, the maximum error is part of the automatic output of the routine, thus requiring no additional computation to obtain it.

3) The problems discussed show clearly the sensitivity
of the method to both the boundary functions and the
geometry of the boundary. The best results seem to
come from boundary functions which are regular in a
large portion of the plane. Poorer results are
obtained for boundary data of low continuity and for
boundary data which possess singularities near or on
the boundary. The method prefers convex regions
(such as the ellipse) over nonconvex regions (such as
the bean), although convexity does not seem to play
as important a role as the smoothness of the boundary
functions.

7.2 Suggestions for Further Research

1) The heart of our computational algorithm is finding
the minimax solution of an overdetermined, inconsistent
system of linear equations. We have presented only
one method of determining this minimax solution,
namely the exchange algorithm. Cheney [3] briefly
discusses two other methods - Polya's algorithm and
the Descent algorithm. A paper* has come to our

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* Bartels, R.H., and Golub, G.H., "Stable Numerical Methods
   for Obtaining the Chebyshev Solution to an Over-
   determined System of Equations", JACM, 11, 1968,
   pp. 401-406.

attention as this thesis is in print which presents
yet another technique. The method uses the exchange
algorithm basically, but computes the inverse anew
each iteration rather than updating the inverse from
the previous iteration. The author claims the method
to be more stable than the exchange algorithm as we
have presented it. Furthermore, the algorithm requires
only that the rank of the coefficient matrix be n,
rather than that the rows of the coefficient matrix
satisfy Haar's condition. A comparative analysis of these
algorithms should be attempted.

2) Certainly, the Dirichlet problem for Laplace's equation,
although important in its own right, is only one problem
in the theory of partial differential equations.
Numerical investigation of the mixed problem and the
Neumann problem for Laplace's equation should be
attempted. Other areas of investigation should include
elliptic equations (of which Laplace's equation is a
subset), boundary-value problems which involve multiply-
connected domains, and the biharmonic equation.

3) We have given no attention to the use for which the
solution of the Dirichlet problem is intended. Many
problems require the integral of the solution over a
region. For these problems, it might be better to
minimize

$$M = \int_C |f(s) - \sum_{j=1}^{n} a_j g_j(s)| ds$$

rather than the maximum error along the boundary.

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APPENDIX

LISTINGS OF FORTRAN IV SUBPROGRAMS

The following pages contain listings of the IBM
System 360 Fortran IV subprograms used in our numerical
experiments.

Subroutines DECOMP, SOLVE, IMPRUV, and INVERT have
been adapted from a set of programs given in Forsythe and
Moler [15]. Details can be found therein. Nevertheless,
for completeness, we give a brief description of each
program.

Subroutine DECOMP (N,A,UL,DET,ISGND) uses Gaussian
elimination to find N-by-N triangular matrices L and U
so that LU = PA , where PA is the matrix A with its
rows interchanged. The matrix L - I + U , where I is
the identity matrix, is stored in UL . DET is the
determinant of A and ISGND is the sign of the determin-
ant of A .

Subroutine SOLVE (N,UL,B,X) uses the LU factorization
from DECOMP to find an approximate solution to a system
of equations AX = B .

Subroutine IMPRUV (N,A,UL,B,X,DIGITS) requires a copy
of the original matrix A , its LU decomposition, a
right-hand side B , and the approximate solution X
computed by SOLVE. It carries out the iterative improve-

ment process, until, if possible, X is nearly accurate to machine precision (about 6 digits on the IBM S/360/67).

Subroutine INVERT (N,A,AINV) uses DECOMP, SOLVE, and IMPRUV to calculate the inverse of A , which is then stored in AINV.

Subroutine MINMAX computes the minimax solution and the associated minimax error of an overdetermined, inconsistent system of linear equations Ax = b .

Subroutine REMES computes the vector $c = (c_1,...,c_n)$ that minimizes

$$M = \max_{(x,y)\epsilon C} \left| f(x,y) - \sum_{j=1}^{n} a_j g_j(x,y) \right| .$$

Subroutine ERRMAX finds $(\bar{x},\bar{y})\epsilon C$ such that

$$\left| f(\bar{x},\bar{y}) - \sum_{j=1}^{n} c_j g_j(\bar{x},\bar{y}) \right| = \max_{(x,y)\epsilon C} \left| f(x,y) - \sum_{j=1}^{n} c_j g_j(x,y) \right| .$$

Subprogram U computes $\sum_{j=1}^{n} c_j g_j(x,y)$.

For the sake of completeness, we have included the sub-programs SGN(X), F(X,Y), G(K,X,Y), and BOUND(X) . SGN(X) is the program for the mathematical function sgn x .

$F(X,Y)$ is the program for the function f specified along the boundary; $G(K,X,Y)$ is the program for the approximating function g_k ; BOUND(X) is the program of the boundary under consideration. These last three programs are valid only for problem 2, function (d) .

Subroutine SING is used by other routines to indicate the occurrence of an error condition.

```fortran
      SUBROUTINE DECOMP(N,A,UL,DET,ISGND)
      DIMENSION A(20,20),UL(20,20),SCALES(20),IPS(20)
      COMMON IPS
      DO 5 I=1,N
      IPS(I)=I
      ROWNRM=0.0
      DO 2 J=1,N
      UL(I,J)=A(I,J)
    2 ROWNRM=AMAX1(ROWNRM,ABS(UL(I,J)))
      IF(ROWNRM)3,4,3
    3 SCALES(I)=1.0/ROWNRM
      GO TO 5
    4 CALL SING(1)
      RETURN
    5 CONTINUE
      DET=1.0
      IS=1
      NM1=N-1
      DO 16 K=1,NM1
      BIG=0.0
      DO 11 I=K,N
      IP=IPS(I)
      SIZE=ABS(UL(IP,K)*SCALES(IP))
      IF(SIZE-BIG)11,11,10
   10 BIG=SIZE
      IDXPIV=I
   11 CONTINUE
      IF(BIG)13,12,13
   12 CALL SING(2)
      RETURN
   13 IF(IDXPIV-K)14,15,14
   14 IS=-IS
      J=IPS(K)
      IPS(K)=IPS(IDXPIV)
      IPS(IDXPIV)=J
   15 KP=IPS(K)
      PIVOT=UL(KP,K)
      DET=DET*PIVOT
      KP1=K+1
      DO 16 I=KP1,N
      IP=IPS(I)
      EM=-UL(IP,K)/PIVOT
      UL(IP,K)=-EM
      DO 16 J=KP1,N
   16 UL(IP,J)=UL(IP,J)+EM*UL(KP,J)
      KP=IPS(N)
      DET=IS*DET*UL(KP,N)
      ISGND=SGN(DET)
      IF(UL(KP,N))19,18,19
   18 CALL SING(2)
```

```
19  RETURN
    END
```

```
      SUBROUTINE SOLVE(N,UL,B,X)
      DIMENSION UL(20,20),B(20),X(20),IPS(20)
      COMMON IPS
      NP1=N+1
      IP=IPS(1)
      X(1)=B(IP)
      DO 2 I=2,N
      IP=IPS(I)
      IM1=I-1
      SUM=0.0
      DO 1 J=1,IM1
1     SUM=SUM+UL(IP,J)*X(J)
2     X(I)=B(IP)-SUM
      IP=IPS(N)
      X(N)=X(N)/UL(IP,N)
      DO 4 IBACK=2,N
      I=NP1-IBACK
      IP=IPS(I)
      IP1=I+1
      SUM=0.0
      DO 3 J=IP1,N
3     SUM=SUM+UL(IP,J)*X(J)
4     X(I)=(X(I)-SUM)/UL(IP,I)
      RETURN
      END
```

```
      SUBROUTINE IMPROV(N,A,UL,B,X,DIGITS)
      DIMENSION A(20,20),UL(20,20),B(20),X(20),R(20),DX(20)
      DOUBLE PRECISION SUM,XJ,AIJ
      EPS=1.0E-6
      ITMAX=12
      XNORM=0.0
      DO 1 I=1,N
    1 XNORM=AMAX1(XNORM,ABS(X(I)))
      IF(XNORM)3,2,3
    2 DIGITS=-ALOG10(EPS)
      GO TO 10
    3 DO 9 ITER=1,ITMAX
      DO 5 I=1,N
      SUM=0.0
      DO 4 J=1,N
      XJ=X(J)
      AIJ=A(I,J)
    4 SUM=SUM+AIJ*XJ
      SUM=B(I)-SUM
    5 R(I)=SUM
      CALL SOLVE(N,UL,R,DX)
      DXNORM=0.0
      DO 6 I=1,N
      T=X(I)
      X(I)=X(I)+DX(I)
      DXNORM=AMAX1(DXNORM,ABS(X(I)-T))
    6 CONTINUE
      IF(ITER-1)8,7,8
    7 DIGITS=-ALOG10(AMAX1(DXNORM/XNORM,EPS))
    8 IF(DXNORM-EPS*XNORM)10,10,9
    9 CONTINUE
      CALL SING(3)
   10 RETURN
      END
```

```
      SUBROUTINE INVERT(N,A,AINV)
      DIMENSION A(20,20),AINV(20,20),B(20),UL(20,20),X(20)
      CALL DECOMP(N,A,UL,DET,ISGND)
      DO 3 J=1,N
      DO 1 I=1,N
      B(I)=0.0
      IF(I.EQ.J)B(I)=1.0
    1 CONTINUE
      CALL SOLVE(N,UL,B,X)
      CALL IMPROV(N,A,UL,B,X,DIGITS)
      DO 2 I=1,N
    2 AINV(I,J)=X(I)
    3 CONTINUE
      RETURN
      END
```

```
      SUBROUTINE MINMAX(M,N,A,B,IV,X)
C
C   M IS THE NUMBER OF EQUATIONS.
C   N IS THE NUMBER OF UNKNOWNS.
C   A IS THE M-BY-N MATRIX OF COEFFICIENTS.
C   B IS THE RIGHT-HAND SIDE.
C   M,N,A,B MUST BE SUPPLIED BY THE CALLING PROGRAM.
C   IV IS THE VECTOR OF INDICES OF THE EQUATIONS
C       PRESENTLY UNDER CONSIDERATION. AN INITIAL
C       GUESS MUST BE SUPPLIED BY THE CALLING
C       PROGRAM. ON RETURN,IV CONTAINS THE INDICES
C       OF THE EQUATIONS OF THE CRITICAL SUBSYSTEM.
C   X IS AN (N+1)-COMPONENT VECTOR.ON RETURN,
C       X(1),...X(N) CONTAIN THE MINIMAX SOLUTION,
C       AND X(N+1) CONTAINS THE MINIMAX ERROR.
C
      DIMENSION A(100,19),B(100),IV(20),X(20)
      DIMENSION AA(20,20),D(20,20),UL(20,20)
      REAL S(20),RHO(100),LAMDA(20)
      DOUBLE PRECISION DKJ,DJK,AIJ,BIK,XJ,LJ,SUM
C
C               STEP 1
C
      NP1=N+1
      DO 10 K=1,NP1
      DO 5 I=1,N
      IF(I-K)1,2,2
    1 IR=IV(I)
      GO TO 3
    2 IR=IV(I+1)
    3 DO 4 J=1,N
    4 AA(I,J)=A(IR,J)
    5 CONTINUE
      CALL DECOMP(N,AA,UL,DET,ISGND)
C
C               STEP 2
C
   10 S(K)=ISGND*(-1)**K
C
C               STEP 3
C
      DO 20 I=1,NP1
      IR=IV(I)
      DO 15 J=1,N
   15 AA(I,J)=A(IR,J)
   20 AA(I,NP1)=S(I)
      CALL INVERT(NP1,AA,D)
C
C               STEP 4
C
```

```
      25 DO 30 K=1,NP1
         SUM=0.0
         DO 29 J=1,NP1
         IR=IV(J)
         DKJ=D(K,J)
         BIR=B(IR)
      29 SUM=SUM+DKJ*BIR
      30 X(K)=SUM
C
C               STEP 5
C
         IALPHA=1
         RHOMAX=0.0
         DO 40 I=1,M
         SUM=0.C
         DO 35 J=1,N
         AIJ=A(I,J)
         XJ=X(J)
      35 SUM=SUM+AIJ*XJ
         RHO(I)=B(I)-SUM
C
C               STEP 6
C
         IF(ABS(RHO(I))-RHOMAX)40,40,39
      39 RHOMAX=ABS(RHO(I))
         IALPHA=I
      40 CONTINUE
C
C               STEP 7
C
         IF(RHOMAX-X(NP1))71,71,42
      42 DO 43 K=1,NP1
         IF(IALPHA-IV(K))43,71,43
      43 CONTINUE
C
C               STEP 8
C
         U=SGN(RHO(IALPHA))
C
C               STEP 9
C
         RMAX=-1E75
         IK=1
         DO 50 K=1,NP1
         SUM=0.0
         DO 45 J=1,N
         DJK=D(J,K)
      45 SUM=SUM+A(IALPHA,J)*DJK
         LAMDA(K)=SUM+U*D(NP1,K)
         IF(D(NP1,K))47,55,47
```

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```
C
C               STEP 10
C
   47 T=LAMDA(K)*U/D(NP1,K)
      IF(T-RMAX)50,50,49
   49 RMAX=T
      IK=K
   50 CONTINUE
      IF(LAMDA(IK))59,55,59
   55 CALL SING(4)
C
C               STEP 11
C
   59 DO 60 K=1,NP1
   60 D(K,IK)=D(K,IK)/DBLE(LAMDA(IK))
      DO 70 K=1,NP1
      DO 70 J=1,NP1
      IF(J-IK)69,70,69
   69 DKJ=D(K,J)
      LJ=LAMDA(J)
      D(K,J)=DKJ-LJ*D(K,IK)
   70 CONTINUE
C
C               STEP 12
C
      IV(IK)=IALPHA
      GO TO 25
   71 RETURN
      END
```

```
      SUBROUTINE REMES(MM,N,PTS,IV,C)
C
C   MM IS THE NUMBER OF POINTS IN THE INITIAL DISCRETE
C        POINT SET.
C   N IS THE NUMBER OF APPROXIMATING FUNCTIONS.
C   PTS IS THE VECTOR OF ABSCISSAS OF THE POINTS IN
C        THE DISCRETE POINT SET.
C   IV IS THE VECTOR OF INDICES OF POINTS IN THE
C        CRITICAL SUBSYSTEM RETURNED BY MINMAX.
C   C IS AN (N+1)-COMPONENT VECTOR.C(1),...C(N)
C        CONTAIN THE MINIMIZING COEFFICIENTS. C(N+1)
C        CONTAINS THE ASSOCIATED MINIMAX ERROR.
C
      DIMENSION PTS(100),IV(20),C(20),A(100,19),B(100)
      M=MM
      EPS=1.0E-3
      ITMAX=90
C
C   EPS AND ITMAX ARE CHOSEN BY THE PROGRAMMER.
C   SET UP THE MATRIX A AND THE VECTOR B FOR MINMAX.
C
      NP1=N+1
      DO 5 I=1,M
      X=PTS(I)
      Y=BOUND(X)
      B(I)=F(X,Y)
      DO 5 J=1,N
    5 A(I,J)=G(J,X,Y)
      DO 25 ITER=1,ITMAX
C
C   SELECT COEFFICIENTS TO MINIMIZE MAXIMUM ERROR
C        ON DISCRETE POINT SET PRESENTLY BEING
C        CONSIDERED.
C
      CALL MINMAX(M,N,A,B,IV,C)
C
C   FIND POINT (CALLED XMAX) OF MAXIMUM ERROR
C        (CALLED EMAX) USING COEFFICIENTS CALCULATED
C        FROM PREVIOUS STATEMENT.
C
      CALL ERRMAX(N,C,XMAX,EMAX)
C
C   EXIT ROUTINE.
C
      IF(EMAX-C(NP1))50,50,15
   15 IF((EMAX-C(NP1))/EMAX-EPS)50,50,16
C
C   INTRODUCE POINT OF MAXIMUM ERROR INTO DISCRETE
C        POINT SET. UPDATE A AND B.
C
```

```
16  M=M+1
    PTS(M)=XMAX
    YMAX=BLLND(XMAX)
    B(M)=F(XMAX,YMAX)
    DO 20 J=1,N
20  A(M,J)=G(J,XMAX,YMAX)
25  CONTINUE
    CALL SING(5)
50  RETURN
    END
```

```
      SUBROUTINE ERRMAX(N,C,XMAX,EMAX)
C
C     A AND B DENOTE THE END POINTS OF THE RANGE OF THE
C        INDEPENDENT VARIABLE. THEY MUST BE SUPPLIED
C        BY THE PROGRAMMER.
C
      DIMENSION C(20)
      PI=3.141593
      N20=20*N
      EMAX=0.0
      DO 10 K=1,N20
C
C     THE KTH ZERO OF THE CHEBYSHEV POLYNOMIAL OF
C        DEGREE 20*N,NORMALIZED TO THE INTERVAL (A,B).
C
      X=0.5*((B-A)*COS(PI*(K-0.5)/N20)+B+A)
      Y=BCOND(X)
      E=ABS(F(X,Y)-G(N,C,X,Y))
      IF(EMAX-E)5,10,10
    5 EMAX=E
      XMAX=X
   10 CONTINUE
      RETURN
      END
```

```
      FUNCTION U(N,CFS,X,Y)
      DIMENSION CFS(20)
      DOUBLE PRECISION SUM,CK
      SUM=0.0
      DO 10 K=1,N
      CK=CFS(K)
   10 SUM=SUM+CK*G(K,X,Y)
      U=SUM
      RETURN
      END
```

```
      FUNCTION SGN(Z)
      IF(Z)1,2,3
1  SGN=-1.0
      RETURN
2  SGN=0.0
      RETURN
3  SGN=1.0
      RETURN
      END
```

```
FUNCTION F(X,Y)
F=ALOG((X+1.5)**2+Y**2)
RETURN
END
```

```
      FUNCTION G(K,X,Y)
      COMPLEX Z
      KM1=K-1
      IF(KM1)1,1,2
    1 G=1.0
      RETURN
    2 Z=CMPLX(X,Y)
      G=REAL(Z**KM1)
      RETURN
      END
```

```
FUNCTION BOUND(X)
BOUND=0.5*SQRT(1-X**2)
RETURN
END
```

```
      SUBROUTINE SING(IWHY)
      GO TO (1,2,3,4,5),IWHY
    1 WRITE(6,11)
      RETURN
    2 WRITE(6,12)
      RETURN
    3 WRITE(6,13)
      RETURN
    4 WRITE(6,14)
      RETURN
    5 WRITE(6,15)
      RETURN
   11 FORMAT('0 MATRIX WITH ZERO ROW IN DECOMP.')
   12 FORMAT('0 SINGULAR MATRIX IN DECOMP.')
   13 FORMAT('0 NO CONVERGENCE IN IMPROV.')
   14 FORMAT('0 HAAR CONDITION VIOLATED IN MINMAX.')
   15 FORMAT('0 NO CONVERGENCE IN REMES.')
      END
```

Minimax Coefficients for the Bean Problem

$$f(x,y) = \begin{cases} 0, & x \le 0 \\ x, & x > 0 \end{cases}$$

n = 3 :	.132231	.461538	.092308		
n = 5 :	.093404	.486649	.008827	.608221	-.034363
n = 7 :	.092597	.536986	-.105155	.653991	.018744
	.180471	.441858			
n = 19:	.075025	.514878	-.209529	.989064	.015350
	-.224194	1.03329	-2.27918	-.035892	.213942
	-2.60345	6.39725	-.026738	1.27328	4.87422
	-9.69454	-.596876	-11.2224	-5.88364	

$$f(x,y) = \begin{cases} 0, & x \le 0 \\ x^3 & x > 0 \end{cases}$$

n = 3 :	.031894	.120151	-.000872		
n = 5 :	.017586	.014766	-.014176	.199221	-.006899
n = 7 :	.013713	.089248	-.032558	.196624	.034272
	.212272	.065763			
n = 19:	.010926	.079131	-.048096	.236794	-.024109
	.241918	.200170	-.175985	.255733	-.415811
	-.305695	.488001	-.453883	1.14545	.212120
	-1.02356	.561388	-1.88455	.951538	

$$f(x,y) = \begin{cases} 0 \,, & x \le 0 \\ x^5, & x > 0 \end{cases}$$

n = 3 :	.007117	.032789	0.0		
n = 5 :	.003361	.023898	-.002281	.056748	.005689
n = 7 :	.002812	.021016	-.007908	.049424	.011138
	.055872	.019046			
n = 19:	.002201	.016757	-.009721	.055003	-.004772
	.072625	.045666	.002280	.071505	-.066786
	-.035986	.078506	-.096603	.184970	.009972
	-.229025	.153825	-.347786	.244229	

$$f(x,y) = e^x \cos y + \log[(1-y)^2 + x^2]$$

n = 3 :	1.19201	.861679	-1.62099		
n = 5 :	1.04665	.993063	-1.80887	1.40812	-.222934
n = 7 :	1.01423	.991961	-1.94136	1.36212	-.017493
	.309424	.610016			
n = 19:	1.00007	.999835	-1.99968	1.49914	-.000545
	.168320	.664599	-.453063	.004362	-.000485
	-.388996	.309699	-.018050	.039101	.243903
	-.168254	.007315	-.114652	-.120182	